## V. DEVANATHAN

# TEXTBOOK OF RELATIVISTIC QUANTUM PHYSICS 



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V. Devanathan

## Textbook of Relativistic Quantum Physics

380 pgs.

## V. Devanathan

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To my esteemed colleague

## Dr. M. Anandakrishnan

who has contributed greatly
to the development of higher education in India

## Preface

Since the publication of my book "Relativistic Quantum Mechanics and Quantum Field Theory" in 2011, two Nobel Prizes have been awarded for the contributions related to the subject matter of this book. One is for the theoretical postulate predicting the Higgs Boson which is detected in the Large Hadron Collider experiments at CERN in the year 2013, almost fifty years after its prediction; and the other is for the experimental discovery of the neutrino oscillations which indicate that the neutrinos are not massless particles but should carry a small but finite mass. The Nobel Prize in Physics for the year 2013 was awarded to Peter Higgs and Francois Englert who proposed in the year 1964, along with four others, the Higgs field and its field quantum - Higgs Boson, after their experimental confirmation by CERN. The Nobel Prize for the year 2015 was awarded to Takaaki Kajita and Arthur B. McDonald, the chief physicists of the two large experimental groups - Super-Kamiokande Collaboration in Japan and the Sudbury Neutrino Observatory in Canada which confirmed without any ambiguity the neutrino oscillation from one flavour to another during its flight. It is interesting to note that the 2013 Nobel Prize is for the confirmation of the Standard Model of elementary particles, which unifies the electro-weak interaction and the strong interactions by the formulation of gauge theories that invoke the Higgs field and Higgs Boson; whereas the Nobel Prize awarded in the year 2015 is for the observation of violation in the Standard Model, which assumes zero rest mass for the neutrinos. The discovery of neutrino oscillations clearly indicate that one should look beyond the Standard Model and attempt to include the gravitational interaction and formulate a unified theory of all the four types of interactions. The String theories and Superstring theories appear to be promising candidates for formulating a theory of everything, of which the Standard Model is the low-energy limit, applicable at energies of the order of 100 GeV .
"Textbook of Relativistic Quantum Physics" is the revised and enlarged edition of the earlier book with the addition of three new chapters. Chapters 1 to 7 deal with Relativistic Quantum Mechanics and chapters 8 to 14 are devoted to Quantum Field Theory.

Special effort is made to highlight Feynman's intuitive approach to Quantum Electrodynamics which has given a visual picture of the processes that we study and a thumb rule for calculating the matrix elements. As Schwinger once remarked, Feynman has indeed carried the Quantum Electrodynamics to the "masses". More than Feynman, it was Dyson who had popularized Feynman diagrams and demonstrated that Feynman's intuitive approach is equivalent to the more general field theoretic approach of Tomonago and Schwinger. You will find in this book the highlights of Feynman's contribution and how it leads to the same results that one can arrive by the Quantum Field Theory.

Let me outline the salient features of this book. Chapter 3 is devoted to the study of the Dirac Equation with external potentials, which enables the analysis of the bound state of hydrogen atom, predicting correctly its spectral lines and their hyperfine splitting. Chapter 4 is on the neutrino and neutrino oscillations. Chapter 13 is on the Gauge Theories which unify the three of the four known elementary particle interactions and lead to the concept of Higgs mechanism by which a particle acquires a mass. The chapters on Neutrino and the Gauge Theories highlight the contributions that led to the award of the Nobel Prizes in the years 2015 and 2013.

The author has not made any serious effort to include all the references. Some original papers have been cited in the text as footnote and some books and review articles are included in the Bibliography at the end.

A notable feature of this book is the inclusion of review questions, problems and solutions to problems at the end of each chapter which will definitely promote a clearer and deeper understanding and appreciation of the subject. The author welcomes any suggestion for improvement and he can be contacted by email: vdevanathan2001@yahoo.com.

One cannot expect the reader to start from the first page and go to the last page in a sequential order. Notations and symbols are defined usually in the text when they are first introduced. The difficulty that a reader usually faces is in deciphering notations and symbols when he wants to refer to any specific topic. To help such a reader, I have appended Appendix A: List of Symbols and Notations.

The author had the benefit of interaction with the students and faculty members of the Department of Nuclear Physics of the University of Madras and also with the participants of the Summer Training Programs in Physics (STPIP) conducted every year jointly under the auspices of the Academy of Sciences, Chennai, the University of Madras, Anna University and the Science City of the Government of Tamil Nadu. Thanks are due to Mr. U. Sagayam, Vice Chairman of the Science City and Dr. J. Devasenapathy for sponsoring the Summer Training Programs. I acknowledge with thanks the cooperation extended to me by my colleagues at the Academy of Sciences Dr. M. Anandakrishnan, Dr. C. Venkateswaran, Dr. P. Maruthamuthu, Dr. R. Asokamani, Dr. A. Stephen, Dr. K. Sivaji and Dr. G. Anbalagan for organizing these summer programs. I am grateful to Dr. M. Balasubramaniam for inviting me to deliver the inaugural lecture at the SERC School on Nuclear Physics held in February, 2017 at Barathiar University; which enabled me to interact with a larger circle of physicists who came from all over India. Useful discussions with Prof. G. Rajasekaran and Prof. D. Indhumathi and the software support extended to me by Mr. S. Gopinath are also acknowledged.

This book is dedicated to my esteemed colleague Dr. M. Anandakrishnan who has contributed greatly to the development of higher education in India. Special thanks are due to Mr. N. K. Mehra, Managing Director of Narosa Publishing House for his interest and cooperation extended by him in the publication of this book.

In conclusion, I wish to thank the members of my family for their tolerance and support that they have extended to me in completing this project.

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## Chapter 1

## Introduction

Classical Mechanics envisages the dichotomy of the universe - matter and waves, mass and energy and space and time. Quantum Mechanics unifies matter and waves whereas Einstein's special theory of relativity unifies mass and energy and also space and time. Relativistic Quantum Mechanics is formulated synthesizing Quantum Mechanics with the special theory of relativity.

Newton's laws of motion and Newtonian mechanics have laid the foundations for classical mechanics but it is found that the Hamiltonian and the Lagrangian formulation of classical mechanics and the Hamilton principle of least action are more suitable for treating complicated problems with constraints. The Hamiltonian theory of classical mechanics has lead to the formulation of Quantum Mechanics by treating the canonically conjugate dynamical variables, position $(x, y, z)$ and momentum $\left(p_{x}, p_{y}, p_{z}\right)$ as operators, obeying certain commutation relations. Using the notation $\left[x, p_{x}\right]_{-}=x p_{x}-p_{x} x$, we have

$$
\begin{equation*}
\left[x, p_{x}\right]_{-}=i \hbar, \quad\left[y, p_{y}\right]_{-}=i \hbar, \quad\left[z, p_{z}\right]_{-}=i \hbar \tag{1.1}
\end{equation*}
$$

Similarly, energy $E$ and time $t$ are conjugate variables, obeying the commutation relation

$$
\begin{equation*}
[E, t]_{-}=i \hbar . \tag{1.2}
\end{equation*}
$$

The commutation relations can be viewed as a consequence of Heisenberg's uncertainty relation that leads to the discrete energy levels of a bound system (i.e. the quantization of energy) and the quantization of angular momentum. This is what is known as the first quantization. The non-relativistic relation between the energy and momentum of a particle
leads to the Schrödinger equation whereas the relativistic relation leads to the Klein-Gordon and Dirac equations.

### 1.1 The Schrödinger Equation

The non-relativistic relation for the energy $E$ of a particle of mass $m$ and momentum $\boldsymbol{p}$ moving in a potential $V$ is given by

$$
\begin{equation*}
\frac{p^{2}}{2 m}+V=E \tag{1.3}
\end{equation*}
$$

Treating $\boldsymbol{p}$ and $E$ as operators

$$
\boldsymbol{p} \longrightarrow-i \hbar \boldsymbol{\nabla} \quad \text { and } \quad E \longrightarrow i \hbar \frac{\partial}{\partial t},
$$

in order to satisfy the commutation relations (1.1) and (1.2), we obtain the Schrödinger wave equation

$$
\begin{equation*}
\left(-\frac{\hbar^{2}}{2 m} \boldsymbol{\nabla}^{2}+V\right) \Psi(\boldsymbol{r}, t)=i \hbar \frac{\partial}{\partial t} \Psi(\boldsymbol{r}, t), \tag{1.4}
\end{equation*}
$$

where $\Psi(\boldsymbol{r}, t)$ is the wave function which is interpreted as the amplitude of probability of finding the particle at position $\boldsymbol{r}$ at time $t$. The complex conjugate of Eq. (1.4) is

$$
\begin{equation*}
\left(-\frac{\hbar^{2}}{2 m} \boldsymbol{\nabla}^{2}+V\right) \Psi^{*}(\boldsymbol{r}, t)=-i \hbar \frac{\partial}{\partial t} \Psi^{*}(\boldsymbol{r}, t) \tag{1.5}
\end{equation*}
$$

Multiplying Eq. (1.4) by $\Psi^{*}$ on the left and Eq. (1.5) by $\Psi$ on the right and subtracting the latter from the former, we get

$$
\begin{align*}
-\frac{\hbar^{2}}{2 m}\left\{\Psi^{*} \nabla^{2} \Psi-\left(\boldsymbol{\nabla}^{2} \Psi^{*}\right) \Psi\right\} & =i \hbar\left(\Psi^{*} \frac{\partial \Psi}{\partial t}+\frac{\partial \Psi^{*}}{\partial t} \Psi\right) ; \\
-\frac{\hbar^{2}}{2 m} \boldsymbol{\nabla} \cdot\left\{\Psi^{*} \boldsymbol{\nabla} \Psi-\left(\boldsymbol{\nabla} \Psi^{*}\right) \Psi\right\} & =i \hbar \frac{\partial}{\partial t}\left(\Psi^{*} \Psi\right) . \tag{1.6}
\end{align*}
$$

Defining the probability current density $\boldsymbol{J}$ and the probability density $\rho$,

$$
\begin{equation*}
\boldsymbol{J}=\frac{\hbar}{2 m i}\left\{\Psi^{*} \nabla \Psi-\left(\nabla \Psi^{*}\right) \Psi\right\} \quad \text { and } \quad \rho=\Psi^{*} \Psi \tag{1.7}
\end{equation*}
$$

we obtain the continuity equation

$$
\begin{equation*}
\frac{\partial \rho}{\partial t}+\boldsymbol{\nabla} \cdot \boldsymbol{J}=0 . \tag{1.8}
\end{equation*}
$$

Equation (1.8) is the continuity equation, expressing the conservation law that the rate of change of particle density in a given region is equivalent to the particle flux through the surface enclosing the region. Please note that the probability density $\rho$ is a positive definite quantity which is physically acceptable.

It is remarkable that the Schrödinger equation can be used to solve a wide range of problems including bound states and scattering states ${ }^{1}$. Only we need to put the appropriate boundary conditions and the appropriate potentials.

### 1.2 The Klein-Gordon Equation

Using the relativistic energy-momentum relation for a free particle with rest mass $m$,

$$
\begin{equation*}
E^{2}=p^{2} c^{2}+m^{2} c^{4}, \tag{1.9}
\end{equation*}
$$

and using the usual recipe of treating the energy and momentum as differential operators

$$
\boldsymbol{p} \longrightarrow-i \hbar \boldsymbol{\nabla} \quad \text { and } \quad E \longrightarrow i \hbar \frac{\partial}{\partial t},
$$

we obtain the relativistic wave equation

$$
\begin{equation*}
-\hbar^{2} \frac{\partial^{2} \Psi}{\partial t^{2}}=\left(-c^{2} \hbar^{2} \nabla^{2}+m^{2} c^{4}\right) \Psi \tag{1.10}
\end{equation*}
$$

which, on rearrangement, yields

$$
\begin{equation*}
\left(\frac{1}{c^{2}} \frac{\partial^{2}}{\partial t^{2}}-\nabla^{2}+\frac{m^{2} c^{2}}{\hbar^{2}}\right) \Psi=0 \tag{1.11}
\end{equation*}
$$

Equation (1.11) is the Klein-Gordon equation, the complex conjugate of which is given by

$$
\begin{equation*}
\left(\frac{1}{c^{2}} \frac{\partial^{2}}{\partial t^{2}}-\nabla^{2}+\frac{m^{2} c^{2}}{\hbar^{2}}\right) \Psi^{*}=0 \tag{1.12}
\end{equation*}
$$

Multiplying Eq. (1.11) by $\Psi^{*}$ on the left and Eq. (1.12) by $\Psi$ on the left and subtracting, we obtain

$$
\begin{equation*}
\frac{1}{c^{2}}\left(\Psi^{*} \frac{\partial^{2} \Psi}{\partial t^{2}}-\Psi \frac{\partial^{2} \Psi^{*}}{\partial t^{2}}\right)-\Psi^{*} \nabla^{2} \Psi+\Psi \nabla^{2} \Psi^{*}=0 \tag{1.13}
\end{equation*}
$$

[^0]Since

$$
\begin{equation*}
\frac{\partial}{\partial t}\left(\Psi^{*} \frac{\partial \Psi}{\partial t}-\Psi \frac{\partial \Psi^{*}}{\partial t}\right)=\Psi^{*} \frac{\partial^{2} \Psi}{\partial t^{2}}-\Psi \frac{\partial^{2} \Psi^{*}}{\partial t^{2}} \tag{1.14}
\end{equation*}
$$

and

$$
\begin{equation*}
\nabla \cdot\left(\Psi^{*} \nabla \Psi-\Psi \nabla \Psi^{*}\right)=\Psi^{*} \nabla^{2} \Psi-\Psi \nabla^{2} \Psi^{*} \tag{1.15}
\end{equation*}
$$

Eq. (1.13) can be written in the form of continuity equation

$$
\begin{equation*}
\frac{\partial \rho}{\partial t}+\boldsymbol{\nabla} \cdot \boldsymbol{J}=0 \tag{1.16}
\end{equation*}
$$

choosing for the current density $\boldsymbol{J}$ the same expression

$$
\begin{equation*}
\boldsymbol{J}=\frac{\hbar}{2 m i}\left\{\Psi^{*} \nabla \Psi-\left(\boldsymbol{\nabla} \Psi^{*}\right) \Psi\right\} \tag{1.1}
\end{equation*}
$$

as given in Eq. (1.7) for the non-relativistic Schrödinger equation. In the present case, the continuity equation will be satisfied only if we choose for the probability density $\rho$ the expression

$$
\begin{equation*}
\rho=\frac{i \hbar}{2 m c^{2}}\left(\Psi^{*} \frac{\partial \Psi}{\partial t}-\Psi \frac{\partial \Psi^{*}}{\partial t}\right) . \tag{1.18}
\end{equation*}
$$

The probability density $\rho$ as given by Eq. (1.18) involves both $\Psi$ and $\partial \Psi / \partial t$ which can be fixed arbitrarily and hence admits both positive and negative values. Since the probability density should be a positivedefinite quantity, the Klein-Gordon equation was not accepted as a wave equation for several years until Pauli and Weisskopf ${ }^{2}$ reinterpreted it as a field equation in the same sense as Maxwell's equation for electromagnetic field. By putting the rest mass $m=0$ in Eq. (1.11), we obtain the field equation for the electromagnetic field.

The Klein-Gordon equation is a second-order differential equation in $t$ and this has yielded the physically unacceptable negative values also for the probability density $\rho$. It may be observed that the Schrödinger equation is a first-order differential equation in $t$ and hence yielded a positivedefinite value for the probability density. Taking this clue, Dirac attempted to linearize the relativistic relation $E^{2}=p^{2} c^{2}+m^{2} c^{4}$ which is quadratic in both $E$ and $\boldsymbol{p}$ and arrived at the Dirac equation.

[^1]Thus the attempts to overcome the early difficulties encountered in the formulation of Relativistic Quantum Mechanics paid rich dividends. Dirac succeeded in linearizing the relativistic relation (1.9) which is quadratic in both energy and momentum and obtained the Dirac equation for the electron with the intrinsic properties of spin and magnetic moment. Interpretation of the Klein-Gordon equation as a field equation has sowed the seed for the development of the Quantum Field theory.

### 1.3 The Dirac Equation

Starting with the relativistic relation $E^{2}=p^{2} c^{2}+m^{2} c^{4}$ between the energy and momentum of a free particle of mass $m$, Dirac obtained a linear relation

$$
\begin{equation*}
E=c \boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m c^{2} \tag{1.19}
\end{equation*}
$$

using two operators $\alpha$ and $\beta$ which commute with both position and momentum vectors. Squaring (1.19), we obtain

$$
\begin{equation*}
E^{2}=c^{2}(\boldsymbol{\alpha} \cdot \boldsymbol{p})^{2}+\beta^{2} m^{2} c^{4}+(\boldsymbol{\alpha} \cdot \boldsymbol{p}) \beta m c^{3}+\beta(\boldsymbol{\alpha} \cdot \boldsymbol{p}) m c^{3}, \tag{1.20}
\end{equation*}
$$

which will reduce to the relativistic energy-momentum relation $E^{2}=$ $p^{2} c^{2}+m^{2} c^{4}$ if the operators $\boldsymbol{\alpha}$ and $\beta$ obey the following relations:

$$
\begin{equation*}
(\boldsymbol{\alpha} \cdot \boldsymbol{p})^{2}=p^{2}, \quad \beta^{2}=1, \quad \boldsymbol{\alpha} \beta=-\beta \boldsymbol{\alpha} \tag{1.21}
\end{equation*}
$$

The relation

$$
(\boldsymbol{\alpha} \cdot \boldsymbol{p})^{2}=\left(\alpha_{x} p_{x}+\alpha_{y} p_{y}+\alpha_{z} p_{z}\right)^{2}=p^{2}
$$

implies

$$
\begin{align*}
& \alpha_{x}^{2}=\alpha_{y}^{2}=\alpha_{z}^{2}=1 ;  \tag{1.22}\\
& \alpha_{x} \alpha_{y}=-\alpha_{y} \alpha_{x} ; \quad \alpha_{y} \alpha_{z}=-\alpha_{z} \alpha_{y} ; \quad \alpha_{z} \alpha_{x}=-\alpha_{x} \alpha_{z} .
\end{align*}
$$

Using the usual recipe for the first quantization of replacing the energy and momentum by differential operators

$$
E \longrightarrow i \hbar \frac{\partial}{\partial t} \quad \text { and } \quad \boldsymbol{p} \longrightarrow-i \hbar \boldsymbol{\nabla}
$$

in Eq. (1.19), the time-dependent Dirac equation for a free particle is obtained.

$$
\begin{equation*}
i \hbar \frac{\partial \Psi(\boldsymbol{r}, t)}{\partial t}=\left(-i c \hbar \boldsymbol{\alpha} \cdot \boldsymbol{\nabla}+\beta m c^{2}\right) \Psi(\boldsymbol{r}, t) . \tag{1.23}
\end{equation*}
$$

The space-time dependence of the free particle Dirac wave function can be explicitly written as

$$
\begin{equation*}
\Psi(\boldsymbol{r}, t)=\psi e^{(i / \hbar)(\boldsymbol{p} \cdot \boldsymbol{r}-E t)} \tag{1.24}
\end{equation*}
$$

This satisfies the time-dependent Dirac equation (1.23) for a free particle; from which we obtain the time-independent Dirac equation

$$
\begin{equation*}
\left(c \boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m c^{2}\right) \psi=E \psi \tag{1.25}
\end{equation*}
$$

with a multi-component wave function $\psi$. If $\alpha_{x}, \alpha_{y}, \alpha_{z}$ and $\beta$ are matrices of dimension $N \times N$, then the Dirac wave function $\psi$ should be a column vector with $N$ components. In Eq. (1.25), $\boldsymbol{p}$ is the momentum vector and not an operator.

## Dirac's $\alpha$ and $\beta$ matrices

The Dirac Hamiltonian should be Hermitian, i.e., $H=H^{\dagger}$. Since $\boldsymbol{p}$ is Hermitian,

$$
H=c(\boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m c) ; \quad H^{\dagger}=c\left(\boldsymbol{\alpha}^{\dagger} \cdot \boldsymbol{p}+\beta^{\dagger} m c\right)
$$

This means that $\alpha_{i},(i=x, y, z)$ and $\beta$ matrices are Hermitian.

$$
\alpha_{i}^{\dagger}=\alpha_{i}, \quad(i=x, y, z) ; \quad \beta^{\dagger}=\beta
$$

Since

$$
\alpha_{i}^{2}=1, \quad(i=x, y, z) \quad \text { and } \beta^{2}=1,
$$

it follows that

$$
\alpha_{i}=\alpha_{i}^{-1} \text { and } \beta=\beta^{-1},
$$

which means that the matrices $\alpha_{i}$ and $\beta$ are non-singular and consequently their determinant is non-zero.

$$
\operatorname{det}\left(\alpha_{i}\right) \neq 0,(i=x, y, z) ; \quad \operatorname{det} \beta \neq 0
$$

We have already seen that the Dirac matrices $\alpha_{x}, \alpha_{y}, \alpha_{z}, \beta$ anticommute. If they are of dimension $N \times N$, then

$$
\begin{align*}
\alpha_{i} \beta & =-\beta \alpha_{i} \\
\operatorname{det}\left(\alpha_{i} \beta\right) & =\operatorname{det}(-1) \operatorname{det}\left(\beta \alpha_{i}\right) \\
\operatorname{det}\left(\alpha_{i}\right) \operatorname{det}(\beta) & =\operatorname{det}(-1) \operatorname{det}(\beta) \operatorname{det}\left(\alpha_{i}\right), \tag{1.26}
\end{align*}
$$

from which, we deduce that

$$
\begin{equation*}
\operatorname{det}(-1)=(-1)^{N}=1, \tag{1.27}
\end{equation*}
$$

assuming that the matrices $\alpha_{i}, \beta$ are of dimensions $N \times N$ and the condition (1.27) is satisfied if $N$ is even.

$$
N=2,4,6, \cdots .
$$

We have seen that the $(2 \times 2)$ matrices permit only three anticommuting matrices (the Pauli matrices) along with the unit matrix as independent matrices and all other matrices can be represented as a linear combination of them. Since in Dirac equation, we have four anticommuting matrices, they should be of higher dimension, at least $(4 \times 4)$.

Dirac's $\alpha$ and $\beta$ matrices are traceless. Since $\alpha_{i} \beta=-\beta \alpha_{i}$, it follows that

$$
\begin{align*}
\alpha_{i} \beta \alpha_{i}^{-1} & =-\beta \\
\text { Trace }\left(\alpha_{i} \beta \alpha_{i}^{-1}\right) & =-\operatorname{Trace} \beta \\
\text { Trace }\left(\alpha_{i}^{-1} \alpha_{i} \beta\right) & =-\operatorname{Trace} \beta \\
\operatorname{Trace}(\beta) & =-\operatorname{Trace} \beta \tag{1.28}
\end{align*}
$$

This means that Trace $\beta=0$. In the above derivation, we have used the cyclic property of the Trace.

$$
\text { Trace }(A B C)=\text { Trace }(C A B)
$$

Similarly, it can be shown that Trace $\alpha_{i}=0,(i=x, y, z)$. Let us now give an explicit representation ${ }^{3}$ for the Dirac matrices which satisfy all the above properties.

$$
\alpha_{i}=\left[\begin{array}{cc}
0 & \sigma_{i}  \tag{1.29}\\
\sigma_{i} & 0
\end{array}\right], \quad \beta=\left[\begin{array}{rr}
I & 0 \\
0 & -I
\end{array}\right],
$$

where $\sigma_{i}$ denotes the familiar Pauli matrices and $I$ denotes the $2 \times 2$ unit matrix. Although the Dirac matrices are of dimension $4 \times 4$, we have

[^2]written conveniently in the $2 \times 2$ form in Eq. (1.29). Explicitly,
\[

$$
\begin{align*}
& \alpha_{x}=\left[\begin{array}{rrrr}
0 & 0 & 0 & 1 \\
0 & 0 & 1 & 0 \\
0 & 1 & 0 & 0 \\
1 & 0 & 0 & 0
\end{array}\right], \quad \alpha_{y}=\left[\begin{array}{rrrr}
0 & 0 & 0 & -i \\
0 & 0 & i & 0 \\
0 & -i & 0 & 0 \\
i & 0 & 0 & 0
\end{array}\right]  \tag{1.30}\\
& \alpha_{z}=\left[\begin{array}{rrrr}
0 & 0 & 1 & 0 \\
0 & 0 & 0 & -1 \\
1 & 0 & 0 & 0 \\
0 & -1 & 0 & 0
\end{array}\right], \quad \beta=\left[\begin{array}{rrrr}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & -1 & 0 \\
0 & 0 & 0 & -1
\end{array}\right] .
\end{align*}
$$
\]

## The continuity equation

Multiplying the Dirac equation (1.23) on the left by $\Psi^{\dagger}$, we get

$$
\begin{equation*}
i \hbar \Psi^{\dagger} \frac{\partial \Psi}{\partial t}=-i c \hbar \Psi^{\dagger} \boldsymbol{\alpha} \cdot \nabla \Psi+\beta m c^{2} \Psi^{\dagger} \Psi \tag{1.31}
\end{equation*}
$$

Taking the Hermitian conjugate of the Dirac equation (1.23) and multiplying on the right by $\Psi$, we get

$$
\begin{equation*}
-i \hbar \frac{\partial \Psi^{\dagger}}{\partial t} \Psi=i c \hbar\left(\nabla \Psi^{\dagger}\right) \cdot \boldsymbol{\alpha}^{\dagger} \Psi+\beta^{\dagger} m c^{2} \Psi^{\dagger} \Psi \tag{1.32}
\end{equation*}
$$

Remembering that $\alpha$ and $\beta$ are Hermitian matrices and subtracting (1.32) from (1.31), we obtain

$$
\begin{align*}
i \hbar \frac{\partial}{\partial t}\left(\Psi^{\dagger} \Psi\right) & =-i c \hbar\left(\Psi^{\dagger} \boldsymbol{\alpha} \cdot \nabla \Psi+\left(\nabla \Psi^{\dagger}\right) \cdot \boldsymbol{\alpha} \Psi\right) \\
& =-i c \hbar \boldsymbol{\nabla} \cdot\left(\Psi^{\dagger} \boldsymbol{\alpha} \Psi\right) \tag{1.33}
\end{align*}
$$

which can be written as a continuity equation

$$
\begin{equation*}
\frac{\partial \rho}{\partial t}+\nabla \cdot \boldsymbol{J}=0 \tag{1.34}
\end{equation*}
$$

with ${ }^{4}$

$$
\begin{equation*}
\rho=\Psi^{\dagger} \Psi=\psi^{\dagger} \psi, \quad \boldsymbol{J}=c \Psi^{\dagger} \boldsymbol{\alpha} \Psi=c \psi^{\dagger} \boldsymbol{\alpha} \psi \tag{1.35}
\end{equation*}
$$

[^3]Since $\psi$ is a column vector and $\psi^{\dagger}$, a row vector as given below

$$
\psi=\left[\begin{array}{l}
a  \tag{1.36}\\
b \\
c \\
d
\end{array}\right], \quad \psi^{\dagger}=\left[\begin{array}{llll}
a^{*} & b^{*} & c^{*} & d^{*}
\end{array}\right]
$$

the probability density

$$
\rho=\Psi^{\dagger} \Psi=\psi^{\dagger} \psi=a^{*} a+b^{*} b+c^{*} c+d^{*} d
$$

is a positive-definite quantity, overcoming the difficulty encountered in the Klein-Gordon equation. The probability current density $\boldsymbol{J}$ together with the probability density $\rho$ obeys the continuity equation (1.34). So, the Dirac equation has become a physically acceptable relativistic wave equation.

## Free particle solutions and their interpretations

The multicomponent wave functions of the Dirac equation are known as Dirac spinors and they correspond to the description of spin $\frac{1}{2} \hbar$ particle. The Dirac Hamiltonian $H=c(\boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m c)$ does not commute with the orbital angular momentum operator $\boldsymbol{L}$ and so the orbital angular momentum $\boldsymbol{L}$ is not a constant of motion but $\boldsymbol{L}+\frac{1}{2} \hbar \boldsymbol{\sigma}$ commutes with $H$ (vide Problem (1.1)). Hence the total angular momentum $\boldsymbol{J}=\boldsymbol{L}+\frac{1}{2} \hbar \boldsymbol{\sigma}$ is a constant of motion, indicating that the Dirac particle has an intrinsic spin $\frac{1}{2} \hbar$. This is the greatest triumph of Dirac equation that it describes a particle such as electron with intrinsic spin $\frac{1}{2} \hbar$ and it is a consequence of the linearization of the relativistic relation $E^{2}=p^{2} c^{2}+m^{2} c^{4}$.

Although the linearization has eliminated the negative values for the probability density and thus yielded physically acceptable solutions of the Dirac equation, yet the Dirac equation permitted both positive and negative energy solutions. The negative energy solutions cannot be thrown away as unphysical since both the positive and negative energy solutions together formed a complete set.

$$
\begin{equation*}
\sum_{i=1}^{i=4} \psi_{i} \psi_{i}^{\dagger}=I \tag{1.37}
\end{equation*}
$$

The Dirac equation has four solutions, of which two of them correspond to positive energies, one with spin up and another with spin down and the
other two correspond to negative energies, one with spin up and another with spin down.

The negative energy solutions do pose a serious problem. If the negative energy states exist, then the electron in a positive energy state will make a transition to the negative energy states and it is impossible to construct any stable atom with electrons. To avoid this disastrous situation, Dirac made a bold suggestion that the negative energy states, though exist, are all filled by electrons and the Dirac vacuum ${ }^{5}$ is one in which all the negative energy states are filled and all the positive energy states are empty. Then the Pauli exclusion principle will prevent the transition of an electron from a positive energy state to the negative energy states which are all occupied and thereby renders stability to the electron in an atom.

## The hole theory of Dirac

What are the consequences of this bold assumption of negative energy sea of filled electrons. A negative energy electron can absorb radiation and jump into a positive energy state; thereby creating a hole in the negative energy sea and a particle in the positive energy state. A hole in the negative energy sea is the absence of a negative energy electron from the Dirac vacuum and hence will be interpreted by the observer as the particle with opposite charge in the positive energy state. Thus the Dirac equation together with the hole theory predicts the existence of positrons which have been subsequently discovered by Anderson ${ }^{6}$. Thus the prediction and the subsequent discovery of positrons support the hole theory of Dirac and introduced the concept of anti-particles, which can be considered as the greatest contribution of Dirac's theory of electron.

A gamma ray of energy greater than $2 m c^{2}$ can excite a negative energy electron to a positive energy state; thereby creating a hole and a particle as illustrated in Fig. 1.1. This process is known as pair creation (creation of positron and electron). The inverse process of pair annihilation (annihilation of positron and electron) accompanied by emission of gamma ray of energy greater than $2 m c^{2}$ is also illustrated in Fig. 1.1.

Although Dirac attempted only to develop a single particle relativistic wave equation, he has ended up with the theory of many particles,

[^4]

Figure 1.1: Electron-positron pair creation and pair annihilation according to Dirac's hole theory
predicting pair creation and pair annihilation. Indeed the Dirac equation together with the hole theory can be truly regarded as the forerunner of quantum field theory.

## Bound state problems

The Dirac equation has emerged as a relativistic wave equation for a spin- $\frac{1}{2}$ particle. Let us use it to find the energy levels of hydrogen atom and also study the properties of neutrino which is a spin- $\frac{1}{2}$ particle with rest mass $m=0$.

The hydrogen atom is a bound state problem wherein the electron which is a spin- $\frac{1}{2}$ particle moves in a central potential. This reduces to the problem of solving the Dirac equation with the Coulomb potential with appropriate boundary conditions. This can be done more easily by using spherical polar coordinates. The solutions yield correctly the energy levels of hydrogen atom with fine structure splitting. This is discussed in Chapter 3.

## The Neutrino

The neutrino is a chargeless massless particle with spin- $\frac{1}{2}$ which was first conjectured in the year 1930 by Wolfgang Pauli in order to preserve the conservation laws of energy, momentum and angular momentum in $\beta$ -
decay. Since it is a spin- $\frac{1}{2}$ particle, the Dirac equation should equally be applicable to the study of neutrino but in this case the rest mass $m=0$. Since neutrinos are chargeless and massless and are produced only in weak interactions and undergo only weak interactions, they are very elusive and cannot be detected directly. They can be detected only indirectly by observing the other particles with which they interact. Neutrinos travel large distances without any interaction and so it took more than 25 years for the experimental confirmation of this particle by Reines and Cowen. They were rewarded with the Nobel Prize in Physics for the year 1956 for this Herculean task.

What is produced in $\beta$-decay is the electron neutrino. Subsequently, it is discovered that there are three flavours of leptons - electron, muon and tau particles with their associated neutrinos. Recently, it is found that the neutrinos, during their long journey without interaction, oscillate from one flavour to another. This is known as neutrino oscillations which has been confirmed by the two large experimental groups - the Super Kamiokande Collaboration in Japan and Sudbury Neutrino Observatory in Canada. The chief physicists of these two groups, Takaaki Kajita and Arther McDonald were rewarded with the Nobel Prize in 2015.

The discovery of neutrino oscillations imply that the neutrinos are not massless particles. They should carry a small mass. There is a new surge of interest in neutrino physics to find experimentally the small mass of the neutrinos and also explore the theoretical consequences. In India, too, efforts are being made to set up an India-based Neutrino Observatory (INO) in Theni District in South India.

### 1.4 Feynman's Positron Theory

We have observed earlier that Dirac equation admits both positive and negative energy solutions and a hole in the negative energy sea is interpreted as positron with positive energy. Feynman developed an unified approach for the treatment of electrons and positrons. One can draw a world line for electron which is propagated from an initial space-time point $\boldsymbol{x}_{1}, t_{1}$ to a final space-time point $\boldsymbol{x}_{2}, t_{2}$ as shown in Fig. 1.2. It is a pictorial representation, known as the Feynman diagram. At any given time, the number of particles need not be conserved but the charge will have to be conserved. Taking this clue, Feynman intuitively identified the positron as the electron travelling backwards in time. Initially, at time $t_{1}$, there is only one electron and finally at time $t_{2}$, there is once again only
one electron but at some intermediate time, $t_{i}$ there are two electrons and one positron. An electron-positron pair has been created at space-time point $\boldsymbol{x}_{A}, t_{A}$, the positron being annihilated by the incoming electron at space-time point $x_{B}, t_{B}$ and the electron emerging into the final state. Feynman describes this scenario beautifully in an inimitable way ${ }^{7}$. "It is as though a bombardier flying low over a road suddenly sees three roads and it is only when two of them come together and disappear again that he realizes that he has simply passed over a long switchback in a single road."


Figure 1.2: Feynman diagram representing a world line for an electron. Between the points, B and A , the electron goes back in time. It is the hole in the negative energy sea which is identified as positron that travels from A to B. At time $t_{1}$, there is only one electron and at time $t_{2}$, there is once again only one electron but at some intermediate time $t_{i}$, there are two electrons and one positron. Please note that the particle number is not conserved but the charge is conserved.

Instead of using the differential equation for the time evolution of the wave function in infinitesimal steps, Feynman developed the concept of propagation kernel which makes a giant leap forward and gives the time evolution of the wave function from one space-time point to another space-time point that is well separated. Feynman's method has greatly simplified the calculation of processes in Quantum Electrodynamics without the need of using the field-theoretical methods. In this book, we shall follow Feynman's method of solving problems in Quantum Electrodynamics and indicate how the Quantum Field Theory and interaction between fields lead to the same results through the S-matrix formalism.

[^5]
### 1.5 Quantum Electrodynamics

In Quantum Electrodynamics (QED), we consider processes involving electrons and positrons, interacting with electromagnetic fields. Rutherford scattering, Compton scattering, electron-electron scattering, electronpositron scattering, electron-positron pair creation, electron-positron pair annihilation are some of the processes considered in this book.

In QED, the interaction between the particles and the electromagnetic fields is regarded as a perturbation. When calculations were done in the lowest order, reasonable results were obtained. Encouraged by the initial success, attempts were made to improve the accuracy by going to higher orders but lo and behold, such attempts encountered infinities! How the infinities are tamed and handled by renormalization procedures forms an interesting anecdote and an important development.

QED was the culmination of the efforts of three theoretical physicists, Sin-Itiro Tomonaga, Julian Schwinger and Richard Feynman who shared the Nobel prize for Physics in the year 1965. Their formulations were so divergent that it was Freeman Dyson ${ }^{8}$ who showed their equivalence. The theory of Feynman differs profoundly in its formulation from that of Tomonaga and Schwinger. The advantages of the Feynman theory are simplicity and ease of applications, while those of Tomonaga-Schwinger are generality and theoretical completeness.

Feynman refers to QED as the jewel of physics - our proudest possession, because QED is able to predict the experimental results with incredible accuracy. They are

1. The Lamb shift in hydrogen spectrum of $2 \mathrm{~S}_{1 / 2}$ level from $2 \mathrm{P}_{1 / 2}$,
2. The anomalous magnetic moment of the electron,
3. Splitting of the ground state of positronium.

What I have referred to above is a success story - success in predicting all the observable quantities with great accuracy by using clever techniques to avoid the infinities. Feynman comments that the renormalization theory is simply a way to sweep difficulties under the rug. All the players, Tomonaga, Schwinger and Feynman feel that the theory that they have developed is intellectually not satisfactory. What they have provided is only a conservative solution but what is needed is a radical innovation and a revolutionary departure similar to what has been made in the nineteen thirties by Bohr, Heisenberg, Schrödinger and Dirac.

[^6]
### 1.6 Quantum Field Theory

Instead of considering the Schrödinger equation, the Klein-Gordon equation and the Dirac equation as single-particle wave equations, they can be treated as classical field equations in the same way as Maxwell's equations for electromagnetic field. These field equations can be deduced using appropriate Lagrangian densities $(\mathscr{L})$ and the Euler-Lagrange equations of motion deduced from the action integral. Since the action integral is relativistically invariant, it is possible to develop a relativistically invariant theory which is manifestly covariant at every stage.

Defining the conjugate momentum field $\Pi(\boldsymbol{x}, t)$

$$
\begin{equation*}
\Pi(\boldsymbol{x}, t)=\frac{\partial \mathscr{L}}{\partial \frac{\partial \Psi}{\partial t}}, \tag{1.38}
\end{equation*}
$$

for each field $\Psi(\boldsymbol{x}, t)$, we can postulate either commutation relations or anticommutation relations for the field operator $\Psi(\boldsymbol{x}, t)$ and its conjugate momentum $\Pi(\boldsymbol{x}, t)$. The following commutation relations

$$
\begin{align*}
& {\left[\Psi(\boldsymbol{x}, t), \Pi\left(\boldsymbol{x}^{\prime}, t\right)\right]_{-}=i \hbar \delta\left(\boldsymbol{x}, \boldsymbol{x}^{\prime}\right),} \\
& {\left[\Psi(\boldsymbol{x}, t), \Psi\left(\boldsymbol{x}^{\prime}, t\right)\right]_{-}=0,}  \tag{1.39}\\
& {\left[\Pi(\boldsymbol{x}, t), \Pi\left(\boldsymbol{x}^{\prime}, t\right)\right]_{-}=0}
\end{align*}
$$

are analogous to the Heisenberg commutation relations used in Quantum Mechanics for $\boldsymbol{x}$ and $\boldsymbol{p}$ and bear a similarity to Poisson brackets in classical mechanics. The commutation relations (1.39) are applicable to the Boson fields whereas the anticommutation relations as shown below have to be used for Fermion fields.

$$
\begin{align*}
& \left\{\Psi(\boldsymbol{x}, t), \Pi\left(\boldsymbol{x}^{\prime}, t\right)\right\}_{+}=i \hbar \delta\left(\boldsymbol{x}, \boldsymbol{x}^{\prime}\right) ; \\
& \left\{\Psi(\boldsymbol{x}, t), \Psi\left(\boldsymbol{x}^{\prime}, t\right)\right\}_{+}=0 ;  \tag{1.40}\\
& \left\{\Pi(\boldsymbol{x}, t), \Pi\left(\boldsymbol{x}^{\prime}, t\right)\right\}_{+}=0 .
\end{align*}
$$

The anticommutation relations have no classical analogue.
The commutation relations (1.39) or anticommutation relations (1.40) are used to quantize the fields and since the quantization is done for the second time, this quantization is known as the second quantization. It may be recalled that the Schrödinger equation, the Klein-Gordon equation and the Dirac equation are themselves obtained earlier from the corresponding Hamiltonians by treating $\boldsymbol{x}$ and $\boldsymbol{p}$ as operators obeying commutation relations. This is what we have referred to as the first quantization.

The field $\Psi(\boldsymbol{x}, t)$ and the conjugate momentum field $\Pi(\boldsymbol{x}, t)$ can be expanded in terms of the annihilation and creation operators (known as Fourier decomposition) and the application of commutation relations (1.39) or anticommutation relations (1.40) yields respectively the commutation relations or anticommutation relations between the annihilation and creation operators.

It is found that the Schrödinger equation, being a non-relativistic one, can be quantized by using either commutation or anticommutation relations. If commutation relations are used, then it becomes applicable to Bosons and if anticommutation relations are used then it becomes applicable to Fermions. But in the case of relativistic field equations, only one of the relations can be used to obtain meaningful results. The Klein-Gordon equation can be quantized only by using the commutation relations between the field function and its corresponding conjugate momentum and is applicable only for Bosons. On the other hand, the Dirac equation can be quantized only by using the anticommutation relations and becomes applicable only to Fermions.

The fields interact with one another and all the phenomena that we observe are due to the interaction of fields. How the interaction of the fields can best be studied through the S -matrix formalism is outlined in Chapter 10.

The S-matrix is a perturbation series. As long as we confine ourselves to the lowest order, we get meaningful results. But, if we try to go to higher orders to improve the accuracy, we face divergent integrals. By using special techniques, known as regularization and renormalization, we can extract finite corrections. But Feynman refers to this as the shell game with the following remarks in his inimitable way.
"The shell game that we play is technically called renormalization. But no matter how clever the word, it is what I would call a dippy process! Having to resort to such hocus-pocus has prevented us from proving that the theory of quantum electrodynamics is mathematically self-consistent. I suspect that renormalization is not mathematically legitimate.... I believe there is really no satisfactory quantum electrodynamics, but I am not sure.... I think that the renormalization theory is simply a way to sweep difficulties of divergences of electrodynamics under the rug. I am, of course not sure of that!"

This scepticism was shared by many at that time. By now, it has become clearer that the renormalization techniques are legitimate methods of arriving at finite results for observable quantities and QED is indeed a
renormalizable field theory. The field theory has been extended to include in its fold weak interactions and strong interactions. This has become possible by the formulation of Gauge theories.

## The Gauge Theories

The Lagrangian densities which serve as the source for all field equations are invariant under certain phase transformations. The invariance under phase transformation leads to certain conservation laws. The phase transformation is often known as the gauge transformation. If the phase transformation is independent of space-time coordinates, it is known as the global gauge transformation. If it depends on the space-time coordinates, then it is known as the local gauge transformation.

Just as the invariance of the Lagrangian under global gauge transformation leads to conserved currents (which is translated into conservation laws), the invariance under local gauge transformation dictates the interaction dynamics. Consider, for example, the Lagrangian for a free electron. The requirement that the Lagrangian should be locally invariant under the same type of transformation can only be fulfilled by introducing additional field, which, in this case turn out to be electromagnetic field. Thus the electromagnetic field is a consequence of the requirement of local symmetry of the Lagrangian for the electron.

This procedure can be generalized to more complex transformations. One starts with a Lagrangian for a matter field and derives the interaction by introducing exactly those fields that make the Lagrangian invariant under a relevant local gauge transformation. Thus, it appears that all fundamental forces can be derived by imposing the condition of local gauge symmetry of matter fields. It is also found that invariance of Lagrangian under gauge transformation is a necessary condition for developing a renormalizable theory.

The gauge transformations can be classified as $\mathrm{U}(1), \mathrm{SU}(2)$ and $\mathrm{SU}(3)$ in the group theory parlance. $\mathrm{U}(1)$ is said to be abelian (commutative) gauge transformation and $\mathrm{SU}(2)$ and $\mathrm{SU}(3)$ are said to be non-abelian (non-commutative) gauge transformation. The invariance under local $\mathrm{U}(1)$ gauge transformation leads to the electromagnetic interaction; the invariance under local $\mathrm{SU}(2) \times \mathrm{U}(1)$ gauge transformation leads to the unified electro-weak interaction and the invariance under local $\mathrm{SU}(3)$ gauge transformation leads to strong interactions (Quantum Chromodynamics). So, the gauge theories successfully explain and unify three out of the four known interactions.

## The Standard Model and beyond

Thus a Standard Model for elementary particle physics has emerged as an acceptable theory for the three of the four known elementary particle interactions. It has been the aim of theoretical physicists to formulate a unified theory which includes also gravitational interaction. It is now being felt that the Quantum Field Theory is a low-energy theory whose validity is limited to 1000 GeV , beyond which one has to look for String Theory with supersymmetry, known as Superstring Theories. This emerging scenario is discussed in the last chapter of this book.

## Review Questions

1.1 What is meant by the first quantization and the second quantization? Distinguish between them and illustrate your answer by giving examples.
1.2 Obtain the Klein-Gordon equation and deduce the continuity equation that it obeys. Show that the Klein-Gordon equation permits the negative values for the probability density.
1.3 Obtain the Dirac equation and show that its solution is a multi-component wave function. Also show that the Dirac equation yields the continuity equation with positive-definite values for the probability density.

## Problems

1.1 Show that the orbital angular momentum operator $\boldsymbol{L}$ does not commute with the Dirac Hamiltonian and only $\boldsymbol{L}+\frac{1}{2} \hbar \boldsymbol{\sigma}$ commutes with the Dirac Hamiltonian; thereby indicating that the Dirac equation is for a particle with intrinsic spin $\frac{1}{2} \hbar \sigma$.
1.2 Show that the electron possesses an intrinsic magnetic moment $M$ due to its spin and that the Dirac equation yields the correct gyromagnetic factor ( $g$ factor) $g=2$.
1.3 Show that $\boldsymbol{\sigma}$ does not commute with the Dirac Hamiltonian but $\boldsymbol{\sigma} \cdot \boldsymbol{p}$ commutes with the Dirac Hamiltonian. Hence demonstrate that the helicity $\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{|p|}$ is a good quantum number for a Dirac particle.

## Solutions to Problems

1.1 It is shown that the orbital angular momentum $(\boldsymbol{L})$ is not a constant of motion since the operator $\boldsymbol{L}$ does not commute with the Dirac Hamiltonian $H$. However, the total angular momentum $\boldsymbol{J}$ which includes the spin $\boldsymbol{S}=\frac{1}{2} \hbar \boldsymbol{\sigma}$ is a constant of motion since $\boldsymbol{J}(=\boldsymbol{L}+\boldsymbol{S})$ commutes with $H$.

The Dirac Hamiltonian and the orbital angular momentum operator are given by

$$
\begin{equation*}
H=c \boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m c^{2}, \quad \boldsymbol{L}=\boldsymbol{r} \times \boldsymbol{p} \tag{1.41}
\end{equation*}
$$

where $\boldsymbol{p}$ is to be regarded as an operator $-i \hbar \boldsymbol{\nabla}$. Since the operator $\boldsymbol{L}$ commutes with the term $\beta m c^{2}$ of the Dirac Hamiltonian,

$$
\begin{align*}
{[H, \boldsymbol{L}]_{-} } & =[c \boldsymbol{\alpha} \cdot \boldsymbol{p}, \boldsymbol{L}]_{-} \\
& =c\left[\boldsymbol{\alpha} \cdot \boldsymbol{p}, L_{x}\right]_{-} \hat{\boldsymbol{e}}_{x}+c\left[\boldsymbol{\alpha} \cdot \boldsymbol{p}, L_{y}\right]_{-} \hat{\boldsymbol{e}}_{y}+c\left[\boldsymbol{\alpha} \cdot \boldsymbol{p}, L_{z}\right]_{-} \hat{\boldsymbol{e}}_{z} \tag{1.42}
\end{align*}
$$

where $\hat{\boldsymbol{e}}_{x}, \hat{\boldsymbol{e}}_{y}, \hat{\boldsymbol{e}}_{z}$ denote the unit vectors. Let us evaluate the first term ${ }^{9}$ on the right hand side of the above equation.

$$
\begin{align*}
c\left[\boldsymbol{\alpha} \cdot \boldsymbol{p}, L_{x}\right]_{-} & =c\left[\alpha_{x} p_{x}+\alpha_{y} p_{y}+\alpha_{z} p_{z}, y p_{z}-z p_{y}\right]_{-} \\
& =c\left[\alpha_{y} p_{y}, y p_{z}\right]_{-}-c\left[\alpha_{z} p_{z}, z p_{y}\right]_{-} \\
& =c \alpha_{y}\left[p_{y}, y\right]_{-} p_{z}-c \alpha_{z}\left[p_{z}, z\right]_{-} p_{y} \\
& =-i \hbar c\left(\alpha_{y} p_{z}-\alpha_{z} p_{y}\right)=-i \hbar c(\boldsymbol{\alpha} \times \boldsymbol{p})_{x} \tag{1.43}
\end{align*}
$$

Similarly, the second and third terms can be evaluated and the final result is

$$
\begin{equation*}
[H, \boldsymbol{L}]_{-}=-i \hbar c(\boldsymbol{\alpha} \times \boldsymbol{p}) \neq 0 \tag{1.44}
\end{equation*}
$$

Let us now evaluate the commutator bracket

$$
\begin{align*}
{[H, \boldsymbol{S}]_{-} } & =\left[H, \frac{1}{2} \hbar \boldsymbol{\sigma}\right]_{-} \\
& =\left[c \boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m c^{2}, \frac{1}{2} \hbar \boldsymbol{\sigma}\right]_{-}=\frac{1}{2} \hbar c[\boldsymbol{\alpha} \cdot \boldsymbol{p}, \boldsymbol{\sigma}]_{-} \tag{1.45}
\end{align*}
$$

Using matrix algebra, we find

$$
(\boldsymbol{\alpha} \cdot \boldsymbol{p}) \boldsymbol{\sigma}=\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \boldsymbol{p}  \tag{1.46}\\
\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0
\end{array}\right]\left[\begin{array}{cc}
\boldsymbol{\sigma} & 0 \\
0 & \boldsymbol{\sigma}
\end{array}\right]=\left[\begin{array}{cc}
0 & (\boldsymbol{\sigma} \cdot \boldsymbol{p}) \boldsymbol{\sigma} \\
(\boldsymbol{\sigma} \cdot \boldsymbol{p}) \boldsymbol{\sigma} & 0
\end{array}\right]
$$

and

$$
\boldsymbol{\sigma}(\boldsymbol{\alpha} \cdot \boldsymbol{p})=\left[\begin{array}{cc}
\boldsymbol{\sigma} & 0  \tag{1.47}\\
0 & \boldsymbol{\sigma}
\end{array}\right]\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \boldsymbol{p} \\
\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0
\end{array}\right]=\left[\begin{array}{cc}
0 & \boldsymbol{\sigma}(\boldsymbol{\sigma} \cdot \boldsymbol{p}) \\
\boldsymbol{\sigma}(\boldsymbol{\sigma} \cdot \boldsymbol{p}) & 0
\end{array}\right]
$$

It can easily seen that

$$
\begin{align*}
(\boldsymbol{\sigma} \cdot \boldsymbol{p}) \boldsymbol{\sigma} & =\boldsymbol{p}+i \boldsymbol{\sigma} \times \boldsymbol{p}  \tag{1.48}\\
\boldsymbol{\sigma}(\boldsymbol{\sigma} \cdot \boldsymbol{p}) & =\boldsymbol{p}-i \boldsymbol{\sigma} \times \boldsymbol{p} \tag{1.49}
\end{align*}
$$

[^7]and consequently
\[

$$
\begin{equation*}
(\boldsymbol{\sigma} \cdot \boldsymbol{p}) \boldsymbol{\sigma}-\boldsymbol{\sigma}(\boldsymbol{\sigma} \cdot \boldsymbol{p})=2 i \boldsymbol{\sigma} \times \boldsymbol{p} \tag{1.50}
\end{equation*}
$$

\]

Substituting the above results, the commutator bracket

$$
\begin{align*}
{[H, \boldsymbol{S}]_{-} } & =\frac{1}{2} \hbar c[\boldsymbol{\alpha} \cdot \boldsymbol{p}, \boldsymbol{\sigma}]_{-} \\
& =\frac{1}{2} \hbar c\left[\begin{array}{cc}
0 & 2 i \boldsymbol{\sigma} \times \boldsymbol{p} \\
2 i \boldsymbol{\sigma} \times \boldsymbol{p} & 0
\end{array}\right] \\
& =i \hbar c(\boldsymbol{\alpha} \times \boldsymbol{p}) . \tag{1.51}
\end{align*}
$$

Thus, we find from Eqs. (1.44) and (1.51) that

$$
\begin{equation*}
[H, \boldsymbol{J}]_{-}=[H, \boldsymbol{L}+\boldsymbol{S}]_{-}=0 . \tag{1.52}
\end{equation*}
$$

This means that the total angular momentum is a constant of motion for the Dirac particle and that the total angular momentum includes the spin of the Dirac particle which is $\frac{1}{2} \hbar$. Thus the Dirac equation is for a particle of spin $\frac{1}{2} \hbar$ and it is a consequence of the linearization of the relativistic energy-momentum relation and writing it as a first order differential equation.
1.2 The Dirac equation for a point particle of charge $e$ in an electromagnetic field, described by the vector and scalar potentials $\boldsymbol{A}$ and $\phi$, is written using the usual recipe $\boldsymbol{p} \rightarrow \boldsymbol{p}-e \boldsymbol{A} / c, E \rightarrow E-e \phi$.

$$
\begin{equation*}
\left\{c \boldsymbol{\alpha} \cdot\left(\boldsymbol{p}-\frac{e \boldsymbol{A}}{c}\right)+\beta m c^{2}\right\} \psi=(E-e \phi) \psi . \tag{1.53}
\end{equation*}
$$

Writing $\psi$ in the two-component form $\left[\begin{array}{l}\varphi_{a} \\ \varphi_{b}\end{array}\right]$ and the matrices $\boldsymbol{\alpha}$ and $\beta$ as $2 \times 2$ matrix, we obtain coupled equations in $\varphi_{a}$ and $\varphi_{b}$.

$$
\begin{array}{r}
\left\{c\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \boldsymbol{P} \\
\boldsymbol{\sigma} \cdot \boldsymbol{P} & 0
\end{array}\right]+\left[\begin{array}{cc}
I & 0 \\
0 & -I
\end{array}\right] m c^{2}\right\}\left[\begin{array}{l}
\varphi_{a} \\
\varphi_{b}
\end{array}\right] \\
 \tag{1.54}\\
=(E-e \phi)\left[\begin{array}{l}
\varphi_{a} \\
\varphi_{b}
\end{array}\right],
\end{array}
$$

using the notation $\boldsymbol{P}=\boldsymbol{p}-e \boldsymbol{A} / c$. The two coupled equations are

$$
\begin{align*}
c(\boldsymbol{\sigma} \cdot \boldsymbol{P}) \varphi_{b}+m c^{2} \varphi_{a} & =(E-e \phi) \varphi_{a}  \tag{1.55}\\
c(\boldsymbol{\sigma} \cdot \boldsymbol{P}) \varphi_{a}-m c^{2} \varphi_{b} & =(E-e \phi) \varphi_{b} \tag{1.56}
\end{align*}
$$

We are interested in the positive energy states with $\varphi_{a}$ as large component and $\varphi_{b}$ as small component. If $\epsilon$ is the kinetic energy of the electron, then

$$
E=m c^{2}+\epsilon
$$

where $m$ is the rest mass of the electron. For small $\epsilon$, we can go to the non-relativistic limit and investigate the Dirac equation in the presence of electromagnetic field. From Eq. (1.56), we get

$$
\begin{equation*}
\varphi_{b}=\frac{c \boldsymbol{\sigma} \cdot \boldsymbol{P}}{E+m c^{2}-e \phi} \varphi_{a} \approx \frac{c \boldsymbol{\sigma} \cdot \boldsymbol{P}}{2 m c^{2}} \varphi_{a}, \quad \text { if } \quad \epsilon, e \phi \ll 2 m c^{2} . \tag{1.57}
\end{equation*}
$$

Substituting this expression for $\varphi_{b}$ in Eq. (1.55), we get

$$
\begin{equation*}
\frac{c^{2}(\boldsymbol{\sigma} \cdot \boldsymbol{P})^{2}}{2 m c^{2}} \varphi_{a}=\left(E-m c^{2}-e \phi\right) \varphi_{a}=(\epsilon-e \phi) \varphi_{a} \tag{1.58}
\end{equation*}
$$

The above equation is the non-relativistic limit of Dirac equation and it is an equation in $\varphi_{a}$. We note that ${ }^{10}$

$$
\begin{align*}
(\boldsymbol{\sigma} \cdot \boldsymbol{P})^{2} & =\boldsymbol{P} \cdot \boldsymbol{P}+i \boldsymbol{\sigma} \cdot(\boldsymbol{P} \times \boldsymbol{P})  \tag{1.59}\\
\boldsymbol{P} \times \boldsymbol{P} & =\left(\boldsymbol{p}-\frac{e \boldsymbol{A}}{c}\right) \times\left(\boldsymbol{p}-\frac{e \boldsymbol{A}}{c}\right) \\
& =-\frac{e}{c}(\boldsymbol{A} \times \boldsymbol{p}+\boldsymbol{p} \times \boldsymbol{A})=i \hbar \frac{e}{c} \boldsymbol{\nabla} \times \boldsymbol{A} . \tag{1.60}
\end{align*}
$$

Substituting (1.59) and (1.60) in Eq. (1.58) and observing that the strength of the magnetic field $\boldsymbol{B}=\boldsymbol{\nabla} \times \boldsymbol{A}$, we get

$$
\begin{equation*}
\left\{\frac{1}{2 m}\left(\boldsymbol{p}-\frac{e \boldsymbol{A}}{c}\right)^{2}-\frac{e \hbar}{2 m c} \boldsymbol{\sigma} \cdot \boldsymbol{B}+e \phi\right\} \varphi_{a}=\epsilon \varphi_{a} . \tag{1.61}
\end{equation*}
$$

One can recognize the second term on the left hand side of the above equation as arising from the magnetic dipole moment

$$
\begin{equation*}
\boldsymbol{M}=\frac{e \hbar}{2 m c} \boldsymbol{\sigma}=\mu_{B} \boldsymbol{\sigma} \tag{1.62}
\end{equation*}
$$

where $\mu_{B}=e \hbar / 2 m c$ is known as the Bohr magneton. Since, for the electron, $e$ is negative, the magnetic moment is antiparallel to the spin. Moreover, since the electron spin is $\boldsymbol{S}=\frac{1}{2} \hbar \boldsymbol{\sigma}$, the magnetic moment is

$$
\begin{equation*}
\boldsymbol{M}=\frac{e}{m c} \boldsymbol{S} \tag{1.63}
\end{equation*}
$$

Thus, we obtain the remarkable result that the gyromagnetic ratio of the electron is $e / m c$ or $g(e / 2 m c)$ with $g=2$. The correct $g$ factor follows directly from the Dirac equation and is not put in by hand from experimental evidence.

[^8]1.3 Let us first evaluate the commutator $[\boldsymbol{\sigma}, H]_{-}$.
\[

$$
\begin{aligned}
{[\boldsymbol{\sigma}, H]_{-} } & =[\boldsymbol{\sigma}, c(\boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m c)]_{-} \\
& =c[\boldsymbol{\sigma}, \boldsymbol{\alpha} \cdot \boldsymbol{p}]_{-}, \quad \text { since } \boldsymbol{\sigma} \text { commutes with } \beta m c \\
& =-2 i c(\boldsymbol{\alpha} \times \boldsymbol{p})
\end{aligned}
$$
\]

using the results (1.46) - (1.50).
In the same way, the commutator $[\boldsymbol{\sigma} \cdot \boldsymbol{p}, H]_{-}$can be evaluated.

$$
\begin{aligned}
{[\boldsymbol{\sigma} \cdot \boldsymbol{p}, H]_{-}=} & {[\boldsymbol{\sigma} \cdot \boldsymbol{p}, c(\boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m c)]_{-} } \\
= & c[\boldsymbol{\sigma} \cdot \boldsymbol{p}, \boldsymbol{\alpha} \cdot \boldsymbol{p}]_{-} \\
= & c\left\{\left[\begin{array}{cc}
\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0 \\
0 & \boldsymbol{\sigma} \cdot \boldsymbol{p}
\end{array}\right]\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \boldsymbol{p} \\
\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0
\end{array}\right]\right. \\
& \left.-\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \boldsymbol{p} \\
\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0
\end{array}\right]\left[\begin{array}{cc}
\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0 \\
0 & \boldsymbol{\sigma} \cdot \boldsymbol{p}
\end{array}\right]\right\} \\
= & c\left\{\left[\begin{array}{cc}
0 & (\boldsymbol{\sigma} \cdot \boldsymbol{p})^{2} \\
(\boldsymbol{\sigma} \cdot \boldsymbol{p})^{2} & 0
\end{array}\right]-\left[\begin{array}{cc}
0 & (\boldsymbol{\sigma} \cdot \boldsymbol{p})^{2} \\
(\boldsymbol{\sigma} \cdot \boldsymbol{p})^{2} & 0
\end{array}\right]\right\} \\
= & 0
\end{aligned}
$$

Since $\boldsymbol{\sigma} \cdot \boldsymbol{p}$ commutes with the Dirac Hamiltonian $H$, one can find simultaneously the eigenvalues of both. Since

$$
\left(\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{|p|}\right)^{2}=1
$$

the eigenvalue of the operator $\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{|p|}$, known as the Helicity operator is

$$
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{|p|}= \pm 1
$$

## Chapter 2

## The Dirac Equation and its Solutions

There are two ways of writing the Dirac equation. One is the conventional way of writing as given by Eq. (2.2), using the Dirac matrices $\alpha$ and $\beta$ and the other is the Feynman way of writing as given by Eq. (2.37), using the $\gamma$ matrices. They differ by a multiplicative factor $\beta\left(\equiv \gamma_{0}\right)$. Their positive energy solutions coincide but their negative energy solutions differ by a change in sign of the momentum vector $\boldsymbol{p}$. In other words, the negative energy solutions of the conventional Dirac equation correspond to electron of momentum $\boldsymbol{p}$ and energy $-E$ whereas the negative energy solutions of Dirac equation in Feynman notation correspond to electron of momentum $-\boldsymbol{p}$ and energy $-E$. The reversal in the sign of momentum for the negative energy electron has been interpreted by Feynman as the reversal of time. Feynman made a bold suggestion that the negative energy electron travels backward in time which is equivalent to the positive energy positron travelling forward in time. This new viewpoint has considerably simplified the calculation of electrodynamic processes which can be elegantly represented by means of Feynman diagrams.

### 2.1 The free particle solutions

Starting with the Dirac Hamiltonian of a free particle with momentum $\boldsymbol{p}$ and rest mass $m$

$$
\begin{equation*}
H=c \boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m c^{2} \tag{2.1}
\end{equation*}
$$

which is linear in energy and momentum, we can write down the timeindependent Dirac equation as

$$
\begin{equation*}
\left(c \boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m c^{2}\right) \psi=E \psi . \tag{2.2}
\end{equation*}
$$

Although $\boldsymbol{\alpha}$ and $\beta$ are $4 \times 4$ matrices as given in (1.30), it is convenient to write them in the compressed form as $2 \times 2$ matrices as shown in Eq. (1.29) for easy algebraic manipulation.

Writing the Dirac equation (2.2) in a more simplified form $A \psi=0$, a non-trivial solution for $\psi$ can be obtained by imposing the condition, $\operatorname{det} A=0$.

$$
\begin{align*}
\operatorname{det} A & =\operatorname{det}\left(c \boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m c^{2}-E I\right) \\
& =\left|\begin{array}{cc}
m c^{2}-E & c \boldsymbol{\sigma} \cdot \boldsymbol{p} \\
c \boldsymbol{\sigma} \cdot \boldsymbol{p} & -\left(m c^{2}+E\right)
\end{array}\right|=0, \tag{2.3}
\end{align*}
$$

where $I$ is the unit matrix and $\boldsymbol{\sigma}\left(\sigma_{x}, \sigma_{y}, \sigma_{z}\right)$ are Pauli matrices. This leads to the relativistic relation $E^{2}=p^{2} c^{2}+m^{2} c^{4}$ and guarantees that the Dirac equation satisfies the relativistic relation. So, the Dirac equation is only a linearized form of the relativistic energy-momentum relation in operator formalism. It follows that the Dirac Hamiltonian has two eigenvalues, one positive and the other negative.

$$
\begin{equation*}
E= \pm \sqrt{p^{2} c^{2}+m^{2} c^{4}}= \pm c W \tag{2.4}
\end{equation*}
$$

with

$$
\begin{equation*}
W=\sqrt{p^{2}+m^{2} c^{2}}=\sqrt{p^{2}+\mu^{2}}, \text { using the symbol } \mu=m c \tag{2.5}
\end{equation*}
$$

In this choice, $W$ is always a positive quantity.
Writing the Dirac equation

$$
\left(c \boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m c^{2}\right) \psi=E \psi
$$

as a $2 \times 2$ matrix equation

$$
\left[\begin{array}{cc}
m c^{2} & c \boldsymbol{\sigma} \cdot \boldsymbol{p}  \tag{2.6}\\
c \boldsymbol{\sigma} \cdot \boldsymbol{p} & -m c^{2}
\end{array}\right]\left[\begin{array}{l}
\varphi_{a} \\
\varphi_{b}
\end{array}\right]=E\left[\begin{array}{c}
\varphi_{a} \\
\varphi_{b}
\end{array}\right],
$$

where we have denoted the Dirac wave function in a two component form $\left[\begin{array}{l}\varphi_{a} \\ \varphi_{b}\end{array}\right]$, we obtain a coupled equation in $\varphi_{a}$ and $\varphi_{b}$.

$$
\begin{align*}
m c^{2} \varphi_{a}+c \boldsymbol{\sigma} \cdot \boldsymbol{p} \varphi_{b} & =E \varphi_{a}  \tag{2.7}\\
c \boldsymbol{\sigma} \cdot \boldsymbol{p} \varphi_{a}-m c^{2} \varphi_{b} & =E \varphi_{b} \tag{2.8}
\end{align*}
$$

The energy $E$ can assume both positive and negative values. The coupled equations (2.7) and (2.8) give only the ratios $\varphi_{a} / \varphi_{b}$. From Eq. (2.7), we obtain

$$
\begin{equation*}
\frac{\varphi_{a}}{\varphi_{b}}=\frac{c \boldsymbol{\sigma} \cdot \boldsymbol{p}}{E-m c^{2}}=\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{ \pm W-\mu}, \tag{2.9}
\end{equation*}
$$

which becomes indeterminate as $\boldsymbol{p} \rightarrow 0$ for positive energy state since both the components $\varphi_{a}$ and $\varphi_{b}$ become zero. So, Eq. (2.9) is relevant only for negative energy states. In a similar way, it can be deduced from Eq. (2.8) that

$$
\begin{equation*}
\frac{\varphi_{b}}{\varphi_{a}}=\frac{c \boldsymbol{\sigma} \cdot \boldsymbol{p}}{E+m c^{2}}=\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{ \pm W+\mu} \tag{2.10}
\end{equation*}
$$

becomes indeterminate for negative energy states as $\boldsymbol{p} \rightarrow 0$ and so applicable only for positive energy states.

So, we shall choose solution (2.10) for positive energy states and solution (2.9) for negative energy states. Since they give only the ratios, one is free to choose either $\varphi_{a}$ or $\varphi_{b}$. The obvious choice is to choose one of them as Pauli spinor,

$$
\chi_{+}=\left[\begin{array}{l}
1 \\
0
\end{array}\right] \quad \text { or } \quad \chi_{-}=\left[\begin{array}{l}
0 \\
1
\end{array}\right] .
$$

For positive energy solutions, let us choose $\varphi_{a}=\chi_{+}$or $\chi_{-}$and for negative energy solutions, we shall choose $\varphi_{b}=\chi_{+}$or $\chi_{-}$. Thus, we obtain the following four solutions:

$$
\begin{align*}
& \psi_{1}=\left[\begin{array}{c}
\chi_{+} \\
\boldsymbol{\sigma} \cdot \boldsymbol{p} \\
W+\mu \\
+
\end{array}\right]=\left[\begin{array}{c}
1 \\
0 \\
\frac{p_{z}}{W+\mu} \\
\frac{p_{+}}{W+\mu}
\end{array}\right] ; \quad \psi_{2}=\left[\begin{array}{c}
\chi_{-} \\
\boldsymbol{\sigma} \cdot \boldsymbol{p} \\
W+\mu \\
-
\end{array}\right]=\left[\begin{array}{c}
0 \\
1 \\
\frac{p_{-}}{W+\mu} \\
\frac{-p_{z}}{W+\mu}
\end{array}\right] ;  \tag{2.11}\\
& \psi_{3}=\left[\begin{array}{c}
\frac{-\boldsymbol{\sigma} \cdot \boldsymbol{p}}{W+\mu} \chi_{+} \\
\chi_{+}
\end{array}\right]=\left[\begin{array}{c}
\frac{-p_{z}}{W+\mu} \\
\frac{-p_{+}}{W+\mu} \\
1 \\
0
\end{array}\right] ; \quad \psi_{4}=\left[\begin{array}{c}
\frac{-\boldsymbol{\sigma} \cdot \boldsymbol{p}}{W+\mu} \chi_{-} \\
\chi_{-}
\end{array}\right]=\left[\begin{array}{c}
\frac{-p_{-}}{W+\mu} \\
\frac{p_{z}}{W+\mu} \\
0 \\
1
\end{array}\right] ; \tag{2.12}
\end{align*}
$$

with

$$
p_{+}=p_{x}+i p_{y} \quad \text { and } \quad p_{-}=p_{x}-i p_{y}
$$

In obtaining the above solutions, we have made use of the following relations:

$$
\begin{align*}
\boldsymbol{\sigma} \cdot \boldsymbol{p} & =\sigma_{x} p_{x}+\sigma_{y} p_{y}+\sigma_{z} p_{z} \\
& =\left[\begin{array}{cc}
p_{z} & p_{x}-i p_{y} \\
p_{x}+i p_{y} & -p_{z}
\end{array}\right]=\left[\begin{array}{cc}
p_{z} & p_{-} \\
p_{+} & -p_{z}
\end{array}\right] .  \tag{2.13}\\
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{W+\mu} \chi_{+} & =\frac{1}{W+\mu}\left[\begin{array}{cc}
p_{z} & p_{-} \\
p_{+} & -p_{z}
\end{array}\right]\left[\begin{array}{l}
1 \\
0
\end{array}\right]=\left[\begin{array}{cc}
\frac{p_{z}}{W_{+}+\mu} \\
\frac{p_{+}}{W+\mu}
\end{array}\right] .  \tag{2.14}\\
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{W+\mu} \chi_{-} & =\frac{1}{W+\mu}\left[\begin{array}{cc}
p_{z} & p_{-} \\
p_{+} & -p_{z}
\end{array}\right]\left[\begin{array}{l}
0 \\
1
\end{array}\right]=\left[\frac{p_{-}}{\frac{W+p_{z}}{W+\mu}}\right] . \tag{2.15}
\end{align*}
$$

The four solutions $\psi_{1}, \psi_{2}, \psi_{3}$ and $\psi_{4}$ of the Dirac equation are called Dirac's four-spinors since the Dirac equation corresponds to spin- $1 / 2$ particle. The solutions $\psi_{1}$ and $\psi_{2}$ correspond to positive energy states and $\psi_{3}$ and $\psi_{4}$ correspond to negative energy states; of which $\psi_{1}$ and $\psi_{3}$ denote the spin-up state and $\psi_{2}$ and $\psi_{4}$ denote the spin-down state.

### 2.2 Orthogonal and closure properties

It can be easily seen that the four solutions are mutually orthogonal. They can also be normalized such that

$$
\left\langle\psi_{i} \mid \psi_{j}\right\rangle=\psi_{i}^{\dagger} \psi_{j}=\delta_{i j}, \quad i, j=1,2,3,4
$$

The normalized solutions of the free particle Dirac equation are

$$
\begin{array}{ll}
\psi_{1}=N\left[\begin{array}{c}
1 \\
0 \\
\frac{p_{z}}{W+\mu} \\
\frac{p_{+}}{W+\mu}
\end{array}\right] ; & \psi_{2}=N\left[\begin{array}{c}
0 \\
1 \\
\frac{p_{-}}{W+\mu} \\
\frac{-p_{z}}{W+\mu}
\end{array}\right] ; \\
\psi_{3}=N\left[\begin{array}{c}
\frac{-z_{z}}{W+\mu} \\
\frac{-p_{+}}{W+\mu} \\
1 \\
0
\end{array}\right] ; & \psi_{4}=N\left[\begin{array}{c}
\frac{-p_{-}}{W+\mu} \\
\frac{p_{z}}{W+\mu} \\
0 \\
1
\end{array}\right] ; \tag{2.17}
\end{array}
$$

with the normalization factor

$$
N=\sqrt{\frac{W+\mu}{2 W}}
$$

Let us now arrange all the four solutions of the Dirac equation for the free particle side by side to form a matrix $D(\boldsymbol{p})$.

$$
D(\boldsymbol{p})=\left(\frac{W+\mu}{2 W}\right)^{1 / 2}\left[\begin{array}{cccc}
1 & 0 & \frac{-p_{z}}{W} & \frac{-p_{-}}{W+\mu}  \tag{2.18}\\
0 & 1 & \frac{-p_{+}}{W+\mu} & \frac{p_{z}}{W+\mu} \\
\frac{p_{z}}{W+\mu} & \frac{p_{-}}{W+\mu} & 1 & 0 \\
\frac{p_{+}}{W+\mu} & \frac{-p_{z}}{W+\mu} & 0 & 1
\end{array}\right] .
$$

If $\psi_{1}(0), \psi_{2}(0), \psi_{3}(0)$ and $\psi_{4}(0)$ are the fundamental four-spinors in the limiting case $\boldsymbol{p} \rightarrow 0$,
$\psi_{1}(0)=\left[\begin{array}{l}1 \\ 0 \\ 0 \\ 0\end{array}\right], \psi_{2}(0)=\left[\begin{array}{l}0 \\ 1 \\ 0 \\ 0\end{array}\right], \quad \psi_{3}(0)=\left[\begin{array}{l}0 \\ 0 \\ 1 \\ 0\end{array}\right], \quad \psi_{4}(0)=\left[\begin{array}{l}0 \\ 0 \\ 0 \\ 1\end{array}\right]$,
then it follows that

$$
\begin{equation*}
\psi_{i}(\boldsymbol{p})=D(\boldsymbol{p}) \psi_{i}(0) \tag{2.20}
\end{equation*}
$$

It can be easily verified that the matrix $D(\boldsymbol{p})$ can be written in a compact form as

$$
\begin{equation*}
D(\boldsymbol{p})=\left(\frac{W+\mu}{2 W}\right)^{1 / 2}\left(1-\frac{\beta(\boldsymbol{\alpha} \cdot \boldsymbol{p})}{W+\mu}\right) \tag{2.21}
\end{equation*}
$$

It is to be observed that $\boldsymbol{p}$ in $D(\boldsymbol{p})$ is not an operator. Equation (2.20) enables you to go from the free particle solutions given in the rest frame to any other frame of reference in which the particle has momentum $\boldsymbol{p}$. So, $D(\boldsymbol{p})$ is known as the Lorentz boost which tells you how the Dirac wave function transforms under Lorentz transformation.

All the four solutions $\psi_{1}, \psi_{2}, \psi_{3}$ and $\psi_{4}$ together form a complete set since

$$
\begin{equation*}
\sum_{i=1}^{4}\left|\psi_{i}\right\rangle\left\langle\psi_{i}\right|=\sum_{i=1}^{4} \psi_{i} \psi_{i}^{\dagger}=I \tag{2.22}
\end{equation*}
$$

This is known as the closure property. That is why the negative energy solutions $\psi_{3}$ and $\psi_{4}$ cannot be discarded as unphysical and the attempt to explain them has led to the important concept of antiparticles.

### 2.3 Projection operators

Instead of summing over all the four states in Eq. (2.22), a partial sum over positive or negative energy states has a nice physical interpretation as projection operator for positive or negative energy states.

$$
\begin{align*}
& \sum_{i=1,2}\left|\psi_{i}\right\rangle\left\langle\psi_{i}\right|=\psi_{1} \psi_{1}^{\dagger}+\psi_{2} \psi_{2}^{\dagger}=\Lambda^{+},  \tag{2.23}\\
& \sum_{i=3,4}\left|\psi_{i}\right\rangle\left\langle\psi_{i}\right|=\psi_{3} \psi_{3}^{\dagger}+\psi_{4} \psi_{4}^{\dagger}=\Lambda^{-}, \tag{2.24}
\end{align*}
$$

where $\Lambda^{+}$and $\Lambda^{-}$are known as projection operators for positive and negative energy states. Substituting the normalized state vectors (2.16) and (2.17) in Eqs. (2.23) and (2.24), we obtain expressions for the projection operators.

$$
\begin{align*}
& \Lambda^{+}=\frac{1}{2}\left(I+\frac{\boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta \mu}{W}\right) .  \tag{2.25}\\
& \Lambda^{-}=\frac{1}{2}\left(I-\frac{\boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta \mu}{W}\right) . \tag{2.26}
\end{align*}
$$

It can be easily verified that $\Lambda^{+}$acting on positive energy states $\psi_{1}$ or $\psi_{2}$ yields unity whereas $\Lambda^{+}$acting on negative energy states $\psi_{3}$ or $\psi_{4}$ yields zero. That is why $\Lambda^{+}$is called the projection operator for positive energy states. An opposite feature is exhibited by $\Lambda^{-}$.

$$
\begin{align*}
& \Lambda^{+} \psi_{i}=\frac{1}{2}\left(I+\frac{\boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta \mu}{W}\right) \psi_{i}= \begin{cases}\psi_{i}, & i=1,2 \\
0, & i=3,4\end{cases}  \tag{2.27}\\
& \Lambda^{-} \psi_{i}=\frac{1}{2}\left(I-\frac{\boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta \mu}{W}\right) \psi_{i}= \begin{cases}0, & i=1,2, \\
\psi_{i}, & i=3,4\end{cases} \tag{2.28}
\end{align*}
$$

It can be easily verified that

$$
\begin{equation*}
\Lambda^{+}+\Lambda^{-}=I \tag{2.29}
\end{equation*}
$$

### 2.4 Sum over spin states

Let us now consider the scattering of a Dirac particle from an initial state $\psi_{i}$ to a final state $\psi_{f}$ due to an interaction represented by the operator $\mathcal{O}$. The transition matrix element $T_{f i}$ is given by ${ }^{1}$

$$
\begin{equation*}
T_{f i}=\left(\bar{\psi}_{f} \mathcal{O} \psi_{i}\right) \tag{2.30}
\end{equation*}
$$

[^9]where $\bar{\psi}_{f}=\psi_{f}^{\dagger} \gamma_{0}$ and $\gamma_{0} \equiv \beta$. If the spins of both the incident and scattered particles are not observed, the transition probability is given by the square of the matrix element, obtained by summing over the final spin states and averaging over the initial spin states.
\[

$$
\begin{align*}
\left|T_{f i}\right|^{2} & =\frac{1}{2} \sum_{i, f}\left(\bar{\psi}_{f} \mathcal{O} \psi_{i}\right)\left(\bar{\psi}_{f} \mathcal{O} \psi_{i}\right)^{\dagger} \\
& =\frac{1}{2} \sum_{i, f}\left(\psi_{f}^{\dagger} \gamma_{0} \mathcal{O} \psi_{i}\right)\left(\psi_{i}^{\dagger} \mathcal{O}^{\dagger} \gamma_{0} \psi_{f}\right), \tag{2.31}
\end{align*}
$$
\]

where the summation indices $i, f$ are over the two spin states denoted by $\psi_{1}$ and $\psi_{2}$, corresponding to positive energy states only. The factor $\frac{1}{2}$ is due to averaging over the two initial spin states $\psi_{1}$ and $\psi_{2}$. Using the algebra of matrix multiplication,

$$
\begin{align*}
\left|T_{f i}\right|^{2} & =\frac{1}{2} \sum_{i, f} \sum_{\rho, \lambda, \rho^{\prime}, \lambda^{\prime}}\left\{\left(\psi_{f}^{\dagger}\right)_{\rho}\left(\gamma_{0} \mathcal{O}\right)_{\rho \lambda}\left(\psi_{i}\right)_{\lambda}\right\}\left\{\left(\psi_{i}^{\dagger}\right)_{\rho^{\prime}}\left(\mathcal{O}^{\dagger} \gamma_{0}\right)_{\rho^{\prime} \lambda^{\prime}}\left(\psi_{f}\right)_{\lambda^{\prime}}\right\} \\
& =\frac{1}{2} \sum_{i, f} \sum_{\rho, \lambda, \rho^{\prime}, \lambda^{\prime}}\left(\psi_{f}^{\dagger}\right)_{\rho}\left(\gamma_{0} \mathcal{O}\right)_{\rho \lambda}\left(\psi_{i} \psi_{i}^{\dagger}\right)_{\lambda \rho^{\prime}}\left(\mathcal{O}^{\dagger} \gamma_{0}\right)_{\rho^{\prime} \lambda^{\prime}}\left(\psi_{f}\right)_{\lambda^{\prime}} \\
& =\frac{1}{2} \sum_{i, f} \sum_{\rho, \lambda, \rho^{\prime}, \lambda^{\prime}}\left(\gamma_{0} \mathcal{O}\right)_{\rho \lambda}\left(\psi_{i} \psi_{i}^{\dagger}\right)_{\lambda \rho^{\prime}}\left(\mathcal{O}^{\dagger} \gamma_{0}\right)_{\rho^{\prime} \lambda^{\prime}}\left(\psi_{f} \psi_{f}^{\dagger}\right)_{\lambda^{\prime} \rho} \\
& =\frac{1}{2} \sum_{\rho, \lambda, \rho^{\prime}, \lambda^{\prime}}\left(\gamma_{0} \mathcal{O}\right)_{\rho \lambda}\left(\Lambda_{i}^{+}\right)_{\lambda \rho^{\prime}}\left(\mathcal{O}^{\dagger} \gamma_{0}\right)_{\rho^{\prime} \lambda^{\prime}}\left(\Lambda_{f}^{+}\right)_{\lambda^{\prime} \rho} \\
& =\frac{1}{2} \operatorname{Tr}\left(\gamma_{0} \mathcal{O} \Lambda_{i}^{+} \mathcal{O}^{\dagger} \gamma_{0} \Lambda_{f}^{+}\right) . \tag{2.32}
\end{align*}
$$

In the above equation, $\Lambda_{i}^{+}$and $\Lambda_{f}^{+}$are the projection operators of the initial and final particle, as defined in Eq. (2.27), obtained after summing over the two spin states corresponding to the positive energy state.

$$
\begin{align*}
\Lambda_{i}^{+} & =\frac{1}{2}\left(I+\frac{\boldsymbol{\alpha} \cdot \boldsymbol{p}_{i}+\beta m}{W_{i}}\right)  \tag{2.33}\\
\Lambda_{f}^{+} & =\frac{1}{2}\left(I+\frac{\boldsymbol{\alpha} \cdot \boldsymbol{p}_{f}+\beta m}{W_{f}}\right) . \tag{2.34}
\end{align*}
$$

### 2.5 In Feynman's notation

Multiply the Dirac equation (2.2) by $\beta$ from the left

$$
\begin{equation*}
\left(c \beta \boldsymbol{\alpha} \cdot \boldsymbol{p}+m c^{2}\right) \psi=\beta E \psi \tag{2.35}
\end{equation*}
$$

Using the natural units $\hbar=1, c=1$ and introducing $\gamma$ matrices

$$
\gamma_{0}=\beta=\left[\begin{array}{rr}
I & 0  \tag{2.36}\\
0 & -I
\end{array}\right], \quad \gamma=\beta \boldsymbol{\alpha}=\left[\begin{array}{rr}
0 & \boldsymbol{\sigma} \\
-\boldsymbol{\sigma} & 0
\end{array}\right],
$$

we can rewrite the Dirac equation after rearrangement in the form

$$
\begin{equation*}
\not p \psi=m \psi, \quad \text { with } \quad \not p=\gamma_{0} E-\gamma \cdot \boldsymbol{p} . \tag{2.37}
\end{equation*}
$$

Writing the characteristic equation for the matrix $\not p$,

$$
\left|\begin{array}{cc}
E-\lambda & -\boldsymbol{\sigma} \cdot \boldsymbol{p}  \tag{2.38}\\
\boldsymbol{\sigma} \cdot \boldsymbol{p} & -E-\lambda
\end{array}\right|=0
$$

we get

$$
-\left(E^{2}-\lambda^{2}\right)+(\boldsymbol{\sigma} \cdot \boldsymbol{p})^{2}=0
$$

from which we obtain the eigenvalues of the matrix $\not p$.

$$
\begin{equation*}
\lambda^{2}=E^{2}-p^{2}=m^{2} ; \quad \lambda= \pm m . \tag{2.39}
\end{equation*}
$$

From (2.39), we find that $E$ can have both positive and negative values.

$$
\begin{equation*}
E^{2}=p^{2}+m^{2} ; \quad E= \pm \sqrt{p^{2}+m^{2}}= \pm E_{p} \tag{2.40}
\end{equation*}
$$

where $E_{p}$ is always a positive quantity. The two eigenvalues of $\not p$ obtained in Eq. (2.39) lead to two equations, one for 'positive eigenvalue' state and the other for 'negative eigenvalue' state.

$$
\begin{align*}
\not p \psi_{p} & =m \psi_{p},  \tag{2.41}\\
\not p \psi_{n} & =-m \psi_{n}, \tag{2.42}
\end{align*}
$$

where $\psi_{p}$ and $\psi_{n}$ denote the positive and negative eigenvalue states. The equation for $\psi_{n}$ is obtained by reversing the sign of energy and momentum so that $\not p$ is changed into $-\not p$. The state $\psi_{n}$ which represents the negative energy electron with momentum $-\boldsymbol{p}$ is to be associated with the state of a positron with positive energy and momentum $+\boldsymbol{p}$, according to the Dirac hole theory.

### 2.5.1 Positive eigenvalue states

Writing Eq. (2.41) in the $2 \times 2$ matrix form, choosing a two component form $\left[\begin{array}{c}\phi_{a} \\ \phi_{b}\end{array}\right]$ for $\psi_{p}$, we get

$$
\left[\begin{array}{cc}
E_{p} & -\boldsymbol{\sigma} \cdot \boldsymbol{p}  \tag{2.43}\\
\boldsymbol{\sigma} \cdot \boldsymbol{p} & -E_{p}
\end{array}\right]\left[\begin{array}{l}
\phi_{a} \\
\phi_{b}
\end{array}\right]=m\left[\begin{array}{c}
\phi_{a} \\
\phi_{b}
\end{array}\right]
$$

which yields two coupled equations in $\phi_{a}$ and $\phi_{b}$.

$$
\begin{align*}
E_{p} \phi_{a}-\boldsymbol{\sigma} \cdot \boldsymbol{p} \phi_{b}=m \phi_{a} & \longrightarrow \quad \frac{\phi_{a}}{\phi_{b}}=\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}-m}  \tag{2.44}\\
\boldsymbol{\sigma} \cdot \boldsymbol{p} \phi_{a}-E_{p} \phi_{b}=m \phi_{b} & \longrightarrow \quad \frac{\phi_{b}}{\phi_{a}}=\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \tag{2.45}
\end{align*}
$$

As $\boldsymbol{p} \rightarrow 0$, Eq. (2.44) becomes indeterminate. So, we choose the solution (2.45) for positive eigenvalue states with $\phi_{a}=\chi_{ \pm}$. Hence

$$
\psi_{p}=\left[\begin{array}{c}
\chi_{ \pm}  \tag{2.46}\\
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{ \pm}
\end{array}\right] .
$$

### 2.5.2 Negative eigenvalue states

In a similar way, let us obtain two coupled equations for $\psi_{n}$ choosing a two component form $\left[\begin{array}{l}\phi_{A} \\ \phi_{B}\end{array}\right]$ for $\psi_{n}$.

$$
\left[\begin{array}{cc}
E_{p} & -\boldsymbol{\sigma} \cdot \boldsymbol{p}  \tag{2.47}\\
\boldsymbol{\sigma} \cdot \boldsymbol{p} & -E_{p}
\end{array}\right]\left[\begin{array}{l}
\phi_{A} \\
\phi_{B}
\end{array}\right]=-m\left[\begin{array}{c}
\phi_{A} \\
\phi_{B}
\end{array}\right],
$$

The coupled equations yield the ratio $\phi_{A} / \phi_{B}$.

$$
\begin{align*}
& E_{p} \phi_{A}-\boldsymbol{\sigma} \cdot \boldsymbol{p} \phi_{B}=-m \phi_{A} \quad \longrightarrow \quad \frac{\phi_{A}}{\phi_{B}}=\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m}  \tag{2.48}\\
& \boldsymbol{\sigma} \cdot \boldsymbol{p} \phi_{A}-E_{p} \phi_{B}=-m \phi_{B} \quad \longrightarrow \quad \frac{\phi_{B}}{\phi_{A}}=\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}-m} \tag{2.49}
\end{align*}
$$

Since (2.49) becomes indeterminate as $\boldsymbol{p} \rightarrow 0$, we shall choose the solution (2.48) with $\phi_{B}=\chi_{ \pm}$. Hence

$$
\psi_{n}=\left[\begin{array}{c}
\frac{\sigma \cdot p}{E_{p}+m} \chi_{ \pm}  \tag{2.50}\\
\chi_{ \pm}
\end{array}\right]
$$

Earlier, we have normalized the Dirac wave functions $\psi^{\dagger} \psi$ to 1 but this normalization is not relativistically invariant. Since $\psi^{\dagger} \psi$ (which is the probability density) transforms as the zeroth component of a four-vector current, it is possible to make a relativistically invariant normalization
by setting it equal to the zeroth component of a suitable four-vector, say, energy-momentum four-vector. Feynman ${ }^{2}$ has chosen the normalization ${ }^{3}$

$$
\begin{equation*}
\psi_{p}^{\dagger} \psi_{p}=2 E_{p} \text { or equivalently } \bar{\psi}_{p} \psi_{p}=\psi_{p}^{\dagger} \gamma_{0} \psi_{p}=2 m \tag{2.51}
\end{equation*}
$$

for positive eigenvalue solutions. The normalized solutions $\psi_{p}$ are

$$
\psi_{p}=\sqrt{E_{p}+m}\left[\begin{array}{c}
\chi_{ \pm}  \tag{2.52}\\
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{ \pm}
\end{array}\right]
$$

It can be easily verified that

$$
\begin{equation*}
\sum_{\text {spins }} \psi_{p} \tilde{\psi}_{p}=\not p+m \tag{2.53}
\end{equation*}
$$

The normalized negative eigenvalue solutions are

$$
\psi_{n}=\sqrt{E_{p}+m}\left[\begin{array}{c}
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{ \pm}  \tag{2.54}\\
\chi_{ \pm}
\end{array}\right]
$$

the normalization being

$$
\begin{equation*}
\psi_{n}^{\dagger} \psi_{n}=2 E_{p} ; \text { or equivalently } \bar{\psi}_{n} \psi_{n}=\psi_{n}^{\dagger} \gamma_{0} \psi_{n}=-2 m . \tag{2.55}
\end{equation*}
$$

It can be easily verified by matrix multiplication that

$$
\begin{equation*}
\sum_{\text {spins }} \psi_{n} \bar{\psi}_{n}=\not p-m \tag{2.56}
\end{equation*}
$$

It can be verified that

$$
\begin{equation*}
\sum_{\text {spins }}\left\{\psi_{p} \bar{\psi}_{p}-\psi_{n} \bar{\psi}_{n}\right\}=2 m I \tag{2.57}
\end{equation*}
$$

The operators $\not p+m$ and $\not p-m$ are the projection operators for the positive and negative eigenvalue states.

$$
\begin{align*}
(\not p+m) \psi_{p} & =2 m \psi_{p}, \\
(p p+m) \psi_{n} & =0,  \tag{2.58}\\
(p p-m) \psi_{p} & =0, \\
(p p-m) \psi_{n} & =-2 m \psi_{n} .
\end{align*}
$$

[^10]Let us now write down explicitly the normalized Feynman wave functions for the two spin orientations ${ }^{4}$ of the positive and negative eigenvalue states.

$$
\begin{array}{ll}
\psi_{1}=\psi_{p \uparrow}=N\left[\begin{array}{c}
1 \\
0 \\
\frac{p_{z}}{E_{p}+m} \\
\frac{p_{+}}{E_{p}+m}
\end{array}\right] ; & \psi_{2}=\psi_{p \downarrow}=N\left[\begin{array}{c}
0 \\
1 \\
\frac{p_{-}}{E_{p}+m} \\
\frac{p_{p}}{E_{p}+m}
\end{array}\right] ; \\
\psi_{3}=\psi_{n \uparrow}=N\left[\begin{array}{c}
\frac{p_{z}}{E_{p}+m} \\
\frac{p_{+}}{E_{p}+m} \\
1 \\
0
\end{array}\right] ; & \psi_{4}=\psi_{n \downarrow}=N\left[\begin{array}{c}
\frac{p_{-}}{E_{p}+m} \\
\frac{p_{z}}{E_{p}+m} \\
0 \\
1
\end{array}\right] ; \tag{2.60}
\end{array}
$$

with the normalization factor

$$
N=\sqrt{E_{p}+m}
$$

Let us now arrange all the four solutions of the Feynman equation for the free particle side by side to form a matrix $S(\boldsymbol{p})$.

$$
S(\boldsymbol{p})=\left(E_{p}+m\right)^{1 / 2}\left[\begin{array}{cccc}
1 & 0 & \frac{p_{z}}{E_{p}+m} & \frac{p_{-}}{E_{p}+m}  \tag{2.61}\\
0 & 1 & \frac{p_{+}}{E_{p}+m} & \frac{p_{p}}{E_{p}+m} \\
\frac{p_{z}}{E_{p}+m} & \frac{p_{-}}{E_{p}+m} & 1 & 0 \\
\frac{p_{+}+m}{E_{p}+m} & \frac{p_{p}}{E_{p}+m} & 0 & 1
\end{array}\right] .
$$

It can be easily verified that the matrix $S(\boldsymbol{p})$ can be written in a compact form as

$$
\begin{equation*}
S(\boldsymbol{p})=\sqrt{E_{p}+m}\left(1+\frac{\boldsymbol{\alpha} \cdot \boldsymbol{p}}{E_{p}+m}\right)=\sqrt{E_{p}+m}\left(1+\frac{\gamma_{0} \boldsymbol{\gamma} \cdot \boldsymbol{p}}{E_{p}+m}\right) . \tag{2.62}
\end{equation*}
$$

If $\psi_{1}(0), \psi_{2}(0), \psi_{3}(0)$ and $\psi_{4}(0)$ are the fundamental four-spinors as given by Eq. (2.19) in the limiting case $\boldsymbol{p} \rightarrow 0$, then

$$
\begin{equation*}
\psi_{i}(\boldsymbol{p})=S(\boldsymbol{p}) \psi_{i}(0) \tag{2.63}
\end{equation*}
$$

Equation (2.63) enables you to go from the free particle solutions given in the rest frame to any other frame of reference in which the particle has momentum $\boldsymbol{p}$. So, $S(\boldsymbol{p})$ is the Lorentz boost operator.

[^11]It may be observed that the positive energy spinors of Dirac given by Eq. (2.16) coincide with the positive eigenvalue spinors of Feynman given by Eq. (2.59), except for normalization factor; whereas the negative energy spinors of Dirac given by Eq. (2.17) differ from the negative eigenvalue spinors of Feynman given by Eq. (2.60) with respect to the sign of the momentum vector $\boldsymbol{p}$. The source of this discrepancy can easily be traced. The negative energy solutions of the Dirac equation are obtained by changing the sign of energy alone and not momentum, whereas in Feynman's negative eigenvalue equation, the signs of both energy and momentum are reversed. A hole in the negative energy sea of electrons with energy $-E_{p}$ and momentum $-\boldsymbol{p}$ is equivalent to a positron with energy $E_{p}$ and momentum $\boldsymbol{p}$. In other words, the absence of an electron with energy $-E_{p}$ and momentum $-\boldsymbol{p}$ in the negative energy sea is interpreted as a positron with energy $E_{p}$ and momentum $\boldsymbol{p}$. If the spin projection of the negative energy electron is $+\frac{1}{2}$, then the associated positron will have spin projection $-\frac{1}{2}$.

For the sake of clarity, we can rewrite the four solutions of the Dirac equations (2.59) and (2.60) in Feynman notation as

$$
\begin{align*}
& \psi_{1}=\psi_{p \uparrow}\left(E_{p}, \boldsymbol{p}\right)=u_{\uparrow}\left(E_{p}, \boldsymbol{p}\right)=u_{1}\left(E_{p}, \boldsymbol{p}\right),  \tag{2.64}\\
& \psi_{2}=\psi_{p \downarrow}\left(E_{p}, \boldsymbol{p}\right)=u_{\downarrow}\left(E_{p}, \boldsymbol{p}\right)=u_{2}\left(E_{p}, \boldsymbol{p}\right),  \tag{2.65}\\
& \psi_{3}=\psi_{n \uparrow}\left(-E_{p},-\boldsymbol{p}\right)=u_{\uparrow}\left(-E_{p},-\boldsymbol{p}\right)=v_{\downarrow}\left(E_{p}, \boldsymbol{p}\right)=v_{2}\left(E_{p}, \boldsymbol{p}\right),  \tag{2.66}\\
& \psi_{4}=\psi_{n \downarrow}\left(-E_{p},-\boldsymbol{p}\right)=u \downarrow\left(-E_{p},-\boldsymbol{p}\right)=v_{\uparrow}\left(E_{p}, \boldsymbol{p}\right)=v_{1}\left(E_{p}, \boldsymbol{p}\right) . \tag{2.67}
\end{align*}
$$

In the above equations, we have introduced a new notation, replacing the positive energy solutions by electron spinors $u_{1}, u_{2}$ and the negative energy solutions by positron spinors $v_{1}, v_{2}$.

Generally, $\psi_{i},(i=1,2,3,4)$ is used to denote the plane wave solutions of the Dirac equation for free particle.

$$
\begin{equation*}
\psi_{i}=u_{i} e^{i \boldsymbol{p} \cdot \boldsymbol{x}} \quad \text { or } \quad \psi_{i}=u_{i} e^{i(\boldsymbol{p} \cdot \boldsymbol{x}-E t)}=u_{i} e^{-i \mathbf{p} \cdot \mathbf{x}} \tag{2.68}
\end{equation*}
$$

depending upon whether we wish to represent the plane wave in the threedimensional space or four-dimensional space. The four dimensional spacetime is preferred for describing any phenomenon in a Lorentz invariant way. Bold italics $\boldsymbol{p}, \boldsymbol{x}$ are used to denote the three-dimensional vectors and bold upright letters $\mathbf{p}, \mathbf{x}$ are used to denote the four-vectors (energymomentum four vector and the space-time four-vector). The plane wave part is suppressed in our discussion till now and we have concentrated only on the spinor part. In the application to Quantum Electrodynamics,
we find it convenient to use the symbols $u_{1}, u_{2}$ to denote electron spinors and $v_{1}, v_{2}$ for positron spinors. The summation over the spins yield

$$
\begin{align*}
& \sum_{\text {spins }} u(\boldsymbol{p}) \bar{u}(\boldsymbol{p})=\sum_{1,2} u_{i}(\boldsymbol{p}) \bar{u}_{i}(\boldsymbol{p})=\not p+m ;  \tag{2.69}\\
& \sum_{\text {spins }} v(\boldsymbol{p}) \bar{v}(\boldsymbol{p})=\sum_{1,2} v_{i}(\boldsymbol{p}) \bar{v}_{i}(\boldsymbol{p})=\not p-m . \tag{2.70}
\end{align*}
$$

### 2.5.3 Sum over spin states

Let us now reconsider the problem of summing over spin states using Feynman's notation. The square of the matrix element (2.31) can be evaluated using Feynman's projection operator for positive energy states.

$$
\begin{align*}
\left|T_{f i}\right|^{2} & =\frac{1}{2} \sum_{i f}\left(\bar{\psi}_{f} \mathcal{O} \psi_{i}\right)\left(\bar{\psi}_{f} \mathcal{O} \psi_{i}\right)^{\dagger} \\
& =\frac{1}{2} \sum_{i f}\left(\bar{\psi}_{f} \mathcal{O} \psi_{i}\right)\left(\psi_{i}^{\dagger} \mathcal{O}^{\dagger} \bar{\psi}_{f}^{\dagger}\right) \\
& =\frac{1}{2} \sum_{i f}\left(\bar{\psi}_{f} \mathcal{O} \psi_{i}\right)\left(\bar{\psi}_{i} \gamma_{0} \mathcal{O}^{\dagger} \gamma_{0} \psi_{f}\right) \\
& =\frac{1}{2} \sum_{i f}\left(\bar{\psi}_{f} \mathcal{O} \psi_{i}\right)\left(\bar{\psi}_{i} \tilde{\mathcal{O}} \psi_{f}\right) \tag{2.71}
\end{align*}
$$

where $\sum_{i f}$ denotes the summation over the positive energy spin states of the incident and scattered particle and $\tilde{\mathcal{O}}$ stands for

$$
\begin{equation*}
\tilde{\mathcal{O}}=\gamma_{0} \mathcal{O}^{\dagger} \gamma_{0} \tag{2.72}
\end{equation*}
$$

Replacing $\sum_{i} \psi_{i} \bar{\psi}_{i}$ by $\left(\not p_{i}+m\right)$ which is the projection operator for positive energy states, we obtain

$$
\begin{equation*}
\left|T_{f i}\right|^{2}=\frac{1}{2} \sum_{f}\left(\bar{\psi}_{f} \mathcal{O}\left(\not p_{i}+m\right) \tilde{\mathcal{O}} \psi_{f}\right) \tag{2.73}
\end{equation*}
$$

Using the algebra of matrix multiplication, we get

$$
\begin{align*}
\left|T_{f i}\right|^{2} & =\frac{1}{2} \sum_{f} \sum_{\rho, \lambda, \rho^{\prime}, \lambda^{\prime}}\left(\bar{\psi}_{f}\right)_{\rho}(\mathcal{O})_{\rho \lambda}\left(\not p_{i}+m\right)_{\lambda \rho^{\prime}}(\tilde{\mathcal{O}})_{\rho^{\prime} \lambda^{\prime}}\left(\psi_{f}\right)_{\lambda^{\prime}} \\
& =\frac{1}{2} \sum_{f} \sum_{\rho, \lambda, \rho^{\prime}, \lambda^{\prime}}(\mathcal{O})_{\rho \lambda}\left(\not p_{i}+m\right)_{\lambda \rho^{\prime}}(\tilde{\mathcal{O}})_{\rho^{\prime} \lambda^{\prime}}\left(\psi_{f} \bar{\psi}_{f}\right)_{\lambda^{\prime} \rho} . \tag{2.74}
\end{align*}
$$

Summing over $f$ over the two positive energy spin states and replacing $\sum_{f} \psi_{f} \bar{\psi}_{f}$ by $\left(\not p_{f}+m\right)$, which is the projection operator for the positive energy states, in Eq. (2.74); we get

$$
\begin{equation*}
\left|T_{f i}\right|^{2}=\frac{1}{2} \operatorname{Tr}\left(\mathcal{O}\left(\not p_{i}+m\right) \tilde{\mathcal{O}}\left(\not p_{f}+m\right)\right) . \tag{2.75}
\end{equation*}
$$

The transition probability per unit time is given by Fermi's golden rule

$$
\begin{equation*}
\text { Transition rate }=2 \pi(\Pi N)^{-1}\left|T_{f i}\right|^{2} \rho_{f}, \tag{2.76}
\end{equation*}
$$

where $\Pi N$ denotes the normalization factor $2 E$ for each of the initial and final particles and $\rho_{f}$ is the density of states for the final particle. The cross section is the transition rate per unit incident flux.

### 2.6 A Consistency Check

We have deduced two different expressions (2.32) and (2.75), for the square of the transition amplitude $\left|T_{f i}\right|^{2}$, one using the Dirac matrices and the other using Feynman's notation. They must be equivalent. To show this, let us start with the projection operator $\Lambda^{+}$for positive energy states as given by Eq. (2.23). (Please note that in natural units $W=E_{p}$.)

$$
\begin{equation*}
\Lambda^{+}=\frac{1}{2}\left(I+\frac{\boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m}{E_{p}}\right)=\frac{1}{2 E_{p}}\left(E_{p}+\boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m\right) . \tag{2.77}
\end{equation*}
$$

Multiply Eq. (2.77) by $\beta^{2}=I$ from the right to obtain

$$
\begin{align*}
\Lambda^{+} & =\frac{1}{2 E_{p}}\left(\beta E_{p}+\boldsymbol{\alpha} \cdot \boldsymbol{p} \beta+m\right) \beta \\
& =\frac{1}{2 E_{p}}\left(\beta E_{p}-\beta \boldsymbol{\alpha} \cdot \boldsymbol{p}+m\right) \beta \\
& =\frac{1}{2 E_{p}}\left(\gamma_{0} E_{p}-\gamma \cdot \boldsymbol{p}+m\right) \beta=\frac{1}{2 E_{p}}(\not p+m) \beta . \tag{2.78}
\end{align*}
$$

Substituting the expression (2.78) for $\Lambda^{+}$, into Eq. (2.32) and remembering that $\beta \equiv \gamma_{0}$ and $\operatorname{Tr}(A B C)=\operatorname{Tr}(B C A)$, we get

$$
\begin{align*}
\left|T_{f i}\right|^{2} & =\frac{1}{8 E_{i} E_{f}} \operatorname{Tr}\left(\gamma_{0} \mathcal{O}\left(\not p_{i}+m\right) \gamma_{0} \mathcal{O}^{\dagger} \gamma_{0}\left(\not p_{f}+m\right) \gamma_{0}\right) \\
& =\frac{1}{8 E_{i} E_{f}} \operatorname{Tr}\left(\mathcal{O}\left(\not p_{i}+m\right) \tilde{\mathcal{O}}\left(\not p_{f}+m\right)\right) \tag{2.79}
\end{align*}
$$

This is identical with Eq. (2.75) except for the additional factor $1 /\left(4 E_{i} E_{f}\right)$ which we include, in Feynman's formalism, as normalization factor $(\Pi N)^{-1}$, as indicated in Eq. (2.76).

### 2.7 Algebra of $\gamma$ Matrices

The square of the transition matrix element given by Eq. (2.75) involves the trace of a product of $\gamma$ matrices. So, it will be fruitful to study the $\gamma$ matrices and their properties for evaluating he traces. The $\gamma$ matrices obey the following relations:

$$
\begin{align*}
& \gamma_{0}^{2}=1, \\
& \gamma_{0} \gamma_{x, y, z}+\gamma_{x, y, z} \gamma_{0}=0,  \tag{2.80}\\
& \gamma_{x} \gamma_{y}+\gamma_{y} \gamma_{x}=0, \quad \gamma_{y}^{2} \gamma_{z}+\gamma_{z}^{2} \gamma_{y}=0, \quad \gamma_{z}^{2}=-1, \\
& \gamma_{z} \gamma_{x}+\gamma_{x} \gamma_{z}=0 .
\end{align*}
$$

Using a unified notation, Eq. (2.80) can be written as

$$
\begin{equation*}
\gamma_{\mu} \gamma_{\nu}+\gamma_{\nu} \gamma_{\mu}=2 g_{\mu \nu}, \tag{2.81}
\end{equation*}
$$

where $g_{\mu \nu}$ is a metric defined by

$$
g_{\mu \nu}=\left\{\begin{align*}
0, & \mu \neq \nu  \tag{2.82}\\
+1, & \mu=\nu=0 \\
-1, & \mu=\nu=x, y, z
\end{align*}\right.
$$

Besides, the matrix $\gamma_{0}$ is Hermitian whereas the matrices $\gamma_{x}, \gamma_{y}, \gamma_{z}$ are anti-Hermitian.

$$
\begin{equation*}
\gamma_{0}^{\dagger}=\gamma_{0} ; \quad \gamma_{k}^{\dagger}=-\gamma_{k}, k=x, y, z . \tag{2.83}
\end{equation*}
$$

It is convenient to define ${ }^{5}$ a matrix $\gamma_{5}$ which occurs frequently.

$$
\gamma_{5}=\gamma_{0} \gamma_{x} \gamma_{y} \gamma_{z}=-i\left[\begin{array}{cc}
0 & I  \tag{2.84}\\
I & 0
\end{array}\right] .
$$

It is easy to verify that

$$
\begin{equation*}
\gamma_{5}^{\dagger}=-\gamma_{5} ; \quad \gamma_{5}^{2}=-1 ; \quad \gamma_{5} \gamma_{\mu}+\gamma_{\mu} \gamma_{5}=0 . \tag{2.85}
\end{equation*}
$$

Following Feynman, we can define $\phi$ as follows:

$$
\begin{equation*}
\not d=a_{0} \gamma_{0}-a_{x} \gamma_{x}-a_{y} \gamma_{y}-a_{z} \gamma_{z} . \tag{2.86}
\end{equation*}
$$

[^12]It can be shown that

$$
\begin{align*}
d \gamma_{5} & =-\gamma_{5} d ;  \tag{2.87}\\
d \not b & =-\not b d+2 \mathbf{a} \cdot \mathbf{b} \quad\left(\mathbf{a} \cdot \mathbf{b}=a_{\mu} b_{\mu}\right) ;  \tag{2.88}\\
\gamma_{x} d \gamma_{x} & =\not d+2 a_{x} \gamma_{x} ;  \tag{2.89}\\
\gamma_{\mu} d \gamma_{\mu} & =-2 \phi ;  \tag{2.90}\\
\gamma_{\mu} d \not b \gamma_{\mu} & =4 \mathbf{a} \cdot \mathbf{b} ;  \tag{2.91}\\
\gamma_{\mu} d \not d \phi \phi \gamma_{\mu} & =-2 \phi \nmid \phi d . \tag{2.92}
\end{align*}
$$

It is important to recall the elementary properties of traces,

$$
\begin{align*}
& \operatorname{Tr}(A B C)=\operatorname{Tr}(B C A)=\operatorname{Tr}(C A B), \\
& \operatorname{Tr}(A+B)=\operatorname{Tr} A+\operatorname{Tr} B, \tag{2.93}
\end{align*}
$$

for evaluating the traces involving a product of $\gamma$ matrices. It is known that the trace of a $\gamma$ matrix is zero.

$$
\begin{equation*}
\operatorname{Tr} \gamma_{\mu}=0 \quad(\mu=0, x, y, z) \tag{2.94}
\end{equation*}
$$

Also the trace of an odd number of $\gamma$ matrices vanishes. To prove this, we start with the relation (2.85) which is equivalent to

$$
\begin{equation*}
\gamma_{5} \gamma_{\mu}\left(\gamma_{5}\right)^{-1}=-\gamma_{\mu} \tag{2.95}
\end{equation*}
$$

It follows that

$$
\begin{equation*}
\gamma_{5} \gamma_{\mu_{1}} \gamma_{\mu_{2}} \cdots \gamma_{\mu_{n}}\left(\gamma_{5}\right)^{-1}=(-1)^{n} \gamma_{\mu_{1}} \gamma_{\mu_{2}} \cdots \gamma_{\mu_{n}} . \tag{2.96}
\end{equation*}
$$

Taking the trace of both sides of Eq. (2.96) and using the elementary property of the trace that $\operatorname{Tr}(A B C)=\operatorname{Tr}(B C A)$, we obtain immediately that

$$
\begin{equation*}
\operatorname{Tr}\left(\gamma_{\mu_{1}} \cdots \gamma_{\mu_{n}}\right)=(-1)^{n} \operatorname{Tr}\left(\gamma_{\mu_{1}} \cdots \gamma_{\mu_{n}}\right) \tag{2.97}
\end{equation*}
$$

Equation (2.97) implies that the trace of an odd number of gamma matrices vanishes.

If $n$ is even, it is always possible to reduce it to $n-2$ factors. For example,

$$
\begin{align*}
\operatorname{Tr}\left(\gamma_{\mu} \gamma_{\nu}\right) & =\operatorname{Tr}\left(\gamma_{\nu} \gamma_{\mu}\right) & & \\
& =\frac{1}{2} \operatorname{Tr}\left(\gamma_{\nu} \gamma_{\mu}+\gamma_{\mu} \gamma_{\nu}\right), & & \text { since } \operatorname{Tr}(\mathrm{AB})=\operatorname{Tr}(\mathrm{BA}) \\
& =g_{\mu \nu} \operatorname{Tr} I, & & \text { using Eq. }(2.81) \\
& =4 g_{\mu \nu} . & & \tag{2.98}
\end{align*}
$$

In a similar way, it can be shown that

$$
\begin{equation*}
\operatorname{Tr}\left(\gamma_{\mu} \gamma_{\nu} \gamma_{\rho} \gamma_{\lambda}\right)=4 g_{\mu \lambda} g_{\nu \rho}-4 g_{\lambda \nu} g_{\rho \mu}+4 g_{\lambda \rho} g_{\mu \nu} \tag{2.99}
\end{equation*}
$$

The following traces which occur frequently are given ${ }^{6}$.

$$
\begin{align*}
\operatorname{Tr}(d \phi b) & \left.=\frac{1}{2} \operatorname{Tr}(\phi \not b+\not b \phi)\right)=\operatorname{Tr}(\mathbf{a} \cdot \mathbf{b})=4 \mathbf{a} \cdot \mathbf{b},  \tag{2.100}\\
\operatorname{Tr}(\nmid \nmid \phi \phi)) & =0 .
\end{align*}
$$

Also, it can be shown that the trace of a product of $\gamma$ matrices is the same as the product of $\gamma$ matrices taken in the reverse order.

$$
\begin{equation*}
\operatorname{Tr}(d \not b \not \subset d)=\operatorname{Tr}(d \phi \nmid \phi d) \tag{2.101}
\end{equation*}
$$

### 2.8 An illustrative example

We shall illustrate the foregoing discussion by evaluating the trace in Eq. (2.75) when the transition operator is given by

$$
\begin{equation*}
\mathcal{O}=\gamma_{0} \tag{2.102}
\end{equation*}
$$

Substituting the operator $\mathcal{O}$ in Eq. (2.75), we obtain

$$
\begin{equation*}
\left|T_{f i}\right|^{2}=\frac{1}{2} \operatorname{Tr}\left(\gamma_{0}\left(\not p_{i}+m\right) \tilde{\gamma}_{0}\left(\not p_{f}+m\right)\right) \tag{2.103}
\end{equation*}
$$

where $\not p=\gamma_{0} E-\boldsymbol{\gamma} \cdot \boldsymbol{p}$. Since $\tilde{\gamma}_{0}=\gamma_{0}$ and the trace of a product of an odd number of $\gamma$ matrices vanishes,

$$
\begin{align*}
\operatorname{Tr}\left(\gamma_{0}\left(\not p_{i}+m\right) \gamma_{0}\left(\not p_{f}+m\right)\right) & =\operatorname{Tr}\left(\gamma_{0} \not p_{i} \gamma_{0} \not p_{f}+m^{2}\right) \\
& =\operatorname{Tr}\left(\left(-\not p_{i} \gamma_{0}+2 E_{i}\right) \gamma_{0} \not p_{f}+m^{2}\right) \\
& =\operatorname{Tr}\left(-\not p_{i} \not p_{f}+2 E_{i} \gamma_{0} \not p_{f}+m^{2}\right) \\
& =-4 \mathbf{p}_{i} \cdot \mathbf{p}_{f}+8 E_{i} E_{f}+4 m^{2} . \tag{2.104}
\end{align*}
$$

Equations (2.98) and (2.100) have been used in deducing the last step in the above equation. Expanding the scalar product of the four-vectors $\mathbf{p}_{i} \cdot \mathbf{p}_{f}=E_{i} E_{f}-\boldsymbol{p}_{i} \cdot \boldsymbol{p}_{f}$ and rearranging, we get

$$
\begin{equation*}
\operatorname{Tr}\left(\gamma_{0}\left(\not p_{i}+m\right) \gamma_{0}\left(\not p_{f}+m\right)\right)=4 E_{i} E_{f}+4 p_{i} p_{f} \cos \theta+4 m^{2} \tag{2.105}
\end{equation*}
$$

where $\theta$ denotes the angle between the two vectors $\boldsymbol{p}_{i}$ and $\boldsymbol{p}_{f}$.

[^13]
## Review Questions

2.1 Write down the Dirac equation for a free particle using the $\alpha$ and $\beta$ matrices and obtain its solutions. How many solutions are there and how are they interpreted? Discuss their orthogonal and closure properties.
2.2 Construct the projection operators for the positive and negative energy states of the Dirac Hamiltonian and demonstrate how they pick up only positive or negative energy states.
2.3 Obtain the Dirac equation in Feynman's notation using $\gamma$ matrices and find its free particle solutions. Show that the positive eigenvalue solutions coincide with those of Dirac whereas the negative eigenvalue solutions differ from the negative energy solutions of Dirac. How will you account for this discrepancy?
2.4 A Dirac particle is scattered from an initial state to a final state due to an interaction represented by the operator $\mathcal{O}$. If the initial and final state spins are not observed, how will you calculate the transition probability by summing over the final spin states and averaging over the initial spin states.
2.5 Write down the gamma matrices and discuss their properties. Show that the trace of a product of an odd number of $\gamma$ matrices vanishes. Evaluate (a) $\operatorname{Tr}\left(\gamma_{\mu} \gamma_{\nu}\right)$ and (b) $\operatorname{Tr}\left(\gamma_{\mu} \gamma_{\nu} \gamma_{\rho} \gamma_{\lambda}\right)$.

## Problems

2.1 Show that the positive eigenvalue states $\psi_{p}$ and the negative eigenvalue states $\psi_{n}$ of the Dirac equation, given by

$$
\psi_{p}=\sqrt{E_{p}+m}\left[\begin{array}{c}
\chi_{ \pm} \\
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E+m} \chi_{ \pm}
\end{array}\right], \quad \psi_{n}=\sqrt{E_{p}+m}\left[\begin{array}{c}
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{ \pm} \\
\chi_{ \pm}
\end{array}\right]
$$

are normalized to yield

$$
\psi_{p}^{\dagger} \psi_{p}=2 E_{p}, \quad \psi_{n}^{\dagger} \psi_{n}=2 E_{p}
$$

and

$$
\bar{\psi}_{p} \psi_{p}=2 m, \quad \bar{\psi}_{n} \psi_{n}=-2 m
$$

2.2 Show that, in Feynman's notation, a sum over spins of positive eigenvalue states and a sum over spins of negative eigenvalue states, done separately, yield

$$
\sum_{\text {spins }} \psi_{p} \bar{\psi}_{p}=\not p+m ; \quad \sum_{\text {spins }} \psi_{n} \bar{\psi}_{n}=\not p-m
$$

Show that they can be considered as unnormalized projection operators for positive and negative eigenvalue states of Dirac equation written in Feynman's notation. Hence obtain the normalized projection operators.
2.3 Given the transition operator $\mathcal{O}=\boldsymbol{\gamma} \cdot \boldsymbol{A}$, where $\boldsymbol{A}$ is a vector in the threedimensional space and not an operator, calculate the transition probability for an electron if the initial and final spin states are not observed.
2.4 Given the transition operator $\mathcal{O}=\gamma_{\mu} J_{\mu}$ where $J_{\mu}$ is a four-vector current, calculate the transition probability for an electron if the initial and final spin states are not observed.
2.5 Show that if the operator $\mathcal{O}=\phi \not b \phi$, then $\tilde{\mathcal{O}}=\gamma_{0} \mathcal{O}^{\dagger} \gamma_{0}=\phi \phi \phi d$.

## Solutions to Problems

2.1 First let us show that the Pauli spin vectors are orthonormal.

$$
\begin{aligned}
& \chi_{+}^{\dagger} \chi_{+}=\left[\begin{array}{ll}
1 & 0
\end{array}\right]\left[\begin{array}{l}
1 \\
0
\end{array}\right]=1 . \quad \chi_{+}^{\dagger} \chi_{-}=\left[\begin{array}{ll}
1 & 0
\end{array}\right]\left[\begin{array}{l}
0 \\
1
\end{array}\right]=0 . \\
& \chi_{-}^{\dagger} \chi_{-}=\left[\begin{array}{ll}
0 & 1
\end{array}\right]\left[\begin{array}{l}
0 \\
1
\end{array}\right]=1 .
\end{aligned}
$$

Altogether, there are four solutions for the Dirac equation. Let us denote the two solutions that correspond to positive eigenvalue, one with spin up and the other with spin down by $\psi_{p \uparrow}, \psi_{p \downarrow}$. The other two solutions that correspond to negative eigenvalue relating to the two spin orientations, spin-up and spin-down are denoted by $\psi_{n \uparrow}, \psi_{n \downarrow}$. Using the orthonormal property of the Pauli spin vectors, it is shown below by direct matrix multiplication that

$$
\begin{aligned}
& \psi_{p \uparrow}^{\dagger} \psi_{p \uparrow}=2 E_{p} . \quad \psi_{n \uparrow}^{\dagger} \psi_{n \uparrow}=2 E_{p} . \\
& \psi_{p \uparrow}^{\dagger} \psi_{p \uparrow}=\left(E_{p}+m\right)\left[\chi_{+}^{\dagger} \frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{+}^{\dagger}\right]\left[\begin{array}{c}
\chi_{+} \\
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \\
+
\end{array}\right] \\
&=\left(E_{p}+m\right)\left\{1+\frac{(\boldsymbol{\sigma} \cdot \boldsymbol{p})^{2}}{\left(E_{p}+m\right)^{2}}\right\} \\
&=\left(E_{p}+m\right)\left\{1+\frac{p^{2}}{\left(E_{p}+m\right)^{2}}\right\} \\
&=\left(E_{p}+m\right)\left\{1+\frac{E_{p}^{2}-m^{2}}{\left(E_{p}+m\right)^{2}}\right\} \\
&= \quad \text { using the relation } E_{p}^{2}=p^{2}+m_{p}^{2} .
\end{aligned}
$$

$$
\begin{aligned}
\psi_{n \uparrow}^{\dagger} \psi_{n \uparrow} & =\left(E_{p}+m\right)\left[\begin{array}{cc}
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{+}^{\dagger} & \chi_{+}^{\dagger}
\end{array}\right]\left[\begin{array}{c}
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{+} \\
\chi_{+}
\end{array}\right] \\
& =\left(E_{p}+m\right)\left\{\frac{(\boldsymbol{\sigma} \cdot \boldsymbol{p})^{2}}{\left(E_{p}+m\right)^{2}}+1\right\} \\
& =\left(E_{p}+m\right)\left\{\frac{p^{2}}{\left(E_{p}+m\right)^{2}}+1\right\} \\
& =2 E_{p}
\end{aligned}
$$

In a similar way, it can be shown that $\psi_{p \downarrow}^{\dagger} \psi_{p \downarrow}=2 E_{p} . \quad \psi_{n \downarrow}^{\dagger} \psi_{n \downarrow}=2 E_{p}$. Further,

$$
\begin{aligned}
& \bar{\psi}_{p \uparrow} \psi_{p \uparrow}=\psi_{p \uparrow}^{\dagger} \gamma_{0} \psi_{p \uparrow} \\
& =\left(E_{p}+m\right)\left[\begin{array}{ll}
\chi_{+}^{\dagger} & \frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{+}^{\dagger}
\end{array}\right]\left[\begin{array}{cc}
I & 0 \\
0 & -I
\end{array}\right]\left[\begin{array}{c}
\chi_{+} \\
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{+}
\end{array}\right] \\
& =\left(E_{p}+m\right)\left[\begin{array}{ll}
\chi_{+}^{\dagger} & \boldsymbol{\sigma} \cdot \boldsymbol{p} \\
E_{p}+m
\end{array} \chi_{+}^{\dagger}\right]\left[\begin{array}{c}
\chi_{+} \\
-\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{+}
\end{array}\right] \\
& =\left(E_{p}+m\right)\left\{1-\frac{(\boldsymbol{\sigma} \cdot \boldsymbol{p})^{2}}{\left(E_{p}+m\right)^{2}}\right\} \\
& =\left(E_{p}+m\right)\left\{1-\frac{p^{2}}{\left(E_{p}+m\right)^{2}}\right\} \\
& =\left(E_{p}+m\right)\left\{1-\frac{E_{p}^{2}-m^{2}}{\left(E_{p}+m\right)^{2}}\right\}, \\
& \text { using the relation } E_{p}^{2}=p^{2}+m^{2} \\
& =2 m \text {. } \\
& \bar{\psi}_{n \uparrow} \psi_{n \uparrow}=\psi_{n \uparrow}^{\dagger} \gamma_{0} \psi_{n \uparrow} \\
& =\left(E_{p}+m\right)\left[\begin{array}{cc}
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{+}^{\dagger} & \chi_{+}^{\dagger}
\end{array}\right]\left[\begin{array}{cc}
I & 0 \\
0 & -I
\end{array}\right]\left[\begin{array}{c}
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{+} \\
\chi_{+}
\end{array}\right] \\
& =\left(E_{p}+m\right)\left[\begin{array}{cc}
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{+}^{\dagger} & \chi_{+}^{\dagger}
\end{array}\right]\left[\begin{array}{c}
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{+} \\
-\chi_{+}
\end{array}\right] \\
& =\left(E_{p}+m\right)\left\{\frac{(\boldsymbol{\sigma} \cdot \boldsymbol{p})^{2}}{\left(E_{p}+m\right)^{2}}-1\right\} \\
& =\left(E_{p}+m\right)\left\{\frac{p^{2}}{\left(E_{p}+m\right)^{2}}-1\right\}=-2 m .
\end{aligned}
$$

Similarly, it follows that $\bar{\psi}_{p \downarrow} \psi_{p \downarrow}=2 m ; \quad \bar{\psi}_{n \downarrow} \psi_{n \downarrow}=-2 m$.
2.2 The quantities $\psi_{p} \bar{\psi}_{p}$ and $\psi_{n} \bar{\psi}_{n}$ are matrices.

It can be easily checked that $\chi_{+} \chi_{+}^{\dagger}+\chi_{-} \chi_{-}^{\dagger}$ is a unit matrix.

$$
\begin{aligned}
& \chi_{+} \chi_{+}^{\dagger}+\chi_{-} \chi_{-}^{\dagger}=\left[\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right] . \\
& \sum_{\text {spins }} \psi_{p} \bar{\psi}_{p}=\left(E_{p}+m\right)\left\{\left[\begin{array}{c}
\overrightarrow{\chi_{+}} \\
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{+}
\end{array}\right]\left[\begin{array}{ll}
\chi_{+}^{\dagger} & \frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{+}^{\dagger}
\end{array}\right]\left[\begin{array}{cc}
I & 0 \\
0 & -I
\end{array}\right]\right. \\
& \left.+\left[\begin{array}{c}
\underset{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{-}
\end{array}\right]\left[\begin{array}{ll}
\chi_{-}^{\dagger} & \frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{-}^{\dagger}
\end{array}\right]\left[\begin{array}{cc}
I & 0 \\
0 & -I
\end{array}\right]\right\} \\
& =\left(E_{p}+m\right)\left[\begin{array}{cc}
1 & \frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \\
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} & \frac{p^{2}}{\left(E_{p}+m\right)^{2}}
\end{array}\right]\left[\begin{array}{cc}
I & 0 \\
0 & -I
\end{array}\right] \\
& =\left(E_{p}+m\right)\left[\begin{array}{cc}
1 & -\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \\
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} & -\frac{p^{2}}{\left(E_{p}+m\right)^{2}}
\end{array}\right] \\
& =\left[\begin{array}{cc}
E_{p}+m & -\boldsymbol{\sigma} \cdot \boldsymbol{p} \\
\boldsymbol{\sigma} \cdot \boldsymbol{p} & -\left(E_{p}-m\right)
\end{array}\right] \\
& =\not p+m,
\end{aligned}
$$

since

$$
\begin{aligned}
\not p+m & =\gamma_{0} E_{p}-\boldsymbol{\gamma} \cdot \boldsymbol{p}+m I \\
& =\left[\begin{array}{cc}
E_{p} & 0 \\
0 & -E_{p}
\end{array}\right]-\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \boldsymbol{p} \\
-\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0
\end{array}\right]+m\left[\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right] .
\end{aligned}
$$

Similarly, it can be shown that

$$
\sum_{\text {spins }} \psi_{n} \bar{\psi}_{n}=\not p-m .
$$

The operators $\not p+m$ and $\not p-m$ are called projection operators since they pick out selectively the positive energy states and negative energy states of the electron.

$$
\begin{array}{ll}
(p p+m) \psi_{p}=2 m \psi_{p} ; & (p p+m) \psi_{n}=0 . \\
(p p-m) \psi_{p}=0 ; & (p p-m) \psi_{n}=-2 m \psi_{n} .
\end{array}
$$

These projection operators are not normalized. The normalized projection operators are $\frac{\not p+m}{2 m}$ for positive energy states and $\frac{m-\not \underline{p}}{2 m}$ for negative energy states.
2.3 From Eq. (2.75), we have

$$
\left|T_{f i}\right|^{2}=\frac{1}{2} \operatorname{Tr}\left\{\mathcal{O}\left(p_{i}+m\right) \tilde{\mathcal{O}}\left(p_{f}+m\right)\right\},
$$

where

$$
\tilde{\mathcal{O}}=\gamma_{0} \mathcal{O}^{\dagger} \gamma_{0}
$$

The transition operator $\mathcal{O}$, in the present case is $\boldsymbol{\gamma} \cdot \boldsymbol{A}$. Substituting it, we shall write down the product of operators $\{\cdots\}$.

$$
\begin{aligned}
\{\cdots\} & =\gamma_{k} A_{k}\left(\gamma_{\mu}\left(p_{i}\right)_{\mu}+m\right) \gamma_{0} \gamma_{l}^{\dagger} A_{l}^{*} \gamma_{0}\left(\gamma_{\nu}\left(p_{f}\right)_{\nu}+m\right) \\
& =\gamma_{k} A_{k}\left(\gamma_{\mu}\left(p_{i}\right)_{\mu}+m\right) \gamma_{l} A_{l}^{*}\left(\gamma_{\nu}\left(p_{f}\right)_{\nu}+m\right) .
\end{aligned}
$$

We have used above the relations $\gamma_{l}^{\dagger}=-\gamma_{l}, \gamma_{0} \gamma_{k}=-\gamma_{k} \gamma_{0}$ and $\gamma_{0}^{2}=1$. Since the trace of a product of odd number of $\gamma$ matrices is zero,

$$
\operatorname{Tr}\{\cdots\}=\operatorname{Tr}\left\{\gamma_{k} \gamma_{\mu} \gamma_{l} \gamma_{\nu} A_{k} A_{l}^{*}\left(p_{i}\right)_{\mu}\left(p_{f}\right)_{\nu}+\gamma_{k} \gamma_{l} A_{k} A_{l}^{*} m^{2}\right\} .
$$

The indices $k$ and $l$ denote the components of a three vector and the indices $\mu$ and $\nu$ denote the components of a four-vector. We have earlier evaluated the traces of even number of $\gamma$ matrices.

$$
\begin{aligned}
& \operatorname{Tr}\left(\gamma_{k} \gamma_{\mu} \gamma_{l} \gamma_{\nu}\right)=4 g_{k \mu} g_{l \nu}-4 g_{k l} g_{\mu \nu}+4 g_{k \nu} g_{\mu l} ; \\
& \operatorname{Tr}\left(\gamma_{k} \gamma_{l}\right)=4 g_{k l} .
\end{aligned}
$$

Using the above results,

$$
\begin{aligned}
\left|T_{f i}\right|^{2}= & \frac{1}{2}\left\{4 g_{k \mu} g_{l \nu}\left(p_{i}\right)_{\mu} A_{k}\left(p_{f}\right)_{\nu} A_{l}^{*}-4 g_{k l} g_{\mu \nu} A_{k} A_{l}^{*}\left(p_{i}\right)_{\mu}\left(p_{f}\right)_{\nu}\right. \\
& \left.+4 g_{k \nu} g_{\mu l}\left(p_{f}\right)_{\nu} A_{k}\left(p_{i}\right)_{\mu} A_{l}^{*}+4 g_{k l} A_{k} A_{l}^{*} m^{2}\right\} \\
= & \frac{1}{2}\left\{4\left(\boldsymbol{p}_{i} \cdot \boldsymbol{A}\right)\left(\boldsymbol{p}_{f} \cdot \boldsymbol{A}^{*}\right)+4\left(\boldsymbol{A} \cdot \boldsymbol{A}^{*}\right)\left(\mathbf{p}_{i} \cdot \mathbf{p}_{f}\right)\right. \\
& \left.+4\left(\boldsymbol{p}_{f} \cdot \boldsymbol{A}\right)\left(\boldsymbol{p}_{i} \cdot \boldsymbol{A}^{*}\right)-4 m^{2} \boldsymbol{A} \cdot \boldsymbol{A}^{*}\right\} \\
= & 2\left\{\left(\boldsymbol{p}_{i} \cdot \boldsymbol{A}\right)\left(\boldsymbol{p}_{f} \cdot \boldsymbol{A}^{*}\right)+\left(\boldsymbol{p}_{f} \cdot \boldsymbol{A}\right)\left(\boldsymbol{p}_{i} \cdot \boldsymbol{A}^{*}\right)\right. \\
& \left.+\left(\boldsymbol{A} \cdot \boldsymbol{A}^{*}\right)\left(\mathbf{p}_{i} \cdot \mathbf{p}_{f}-m^{2}\right)\right\},
\end{aligned}
$$

where $\mathbf{p}_{i} \cdot \mathbf{p}_{f}$ denotes the scalar product of four vectors, defined by

$$
\mathbf{p}_{i} \cdot \mathbf{p}_{f}=E_{i} E_{f}-\boldsymbol{p}_{i} \cdot \boldsymbol{p}_{f}
$$

2.4 From Eq. (2.75), we have

$$
\left|T_{f i}\right|^{2}=\frac{1}{2} \operatorname{Tr}\left\{\mathcal{O}\left(p_{i}+m\right) \tilde{\mathcal{O}}\left(p_{f}+m\right)\right\},
$$

where

$$
\tilde{\mathcal{O}}=\gamma_{0} \mathcal{O}^{\dagger} \gamma_{0}
$$

The transition operator $\mathcal{O}$, in the present case is $\gamma_{\mu} J_{\mu}$ where $J_{\mu}$ is a fourvector current.

$$
\tilde{\mathcal{O}}=\gamma_{0} \gamma_{\mu}^{\dagger} \gamma_{0} J_{\mu}^{*}=\gamma_{\mu} J_{\mu}^{*}
$$

Substituting them, we shall write down the product of operators $\{\cdots\}$.

$$
\{\cdots\}=\left[\gamma_{\mu}\left(\not p_{i}+m\right) \gamma_{\nu}\left(p_{f}+m\right)\right] J_{\mu} J_{\nu}^{*}
$$

Since the trace of a product of odd number of $\gamma$ matrices is zero, we find

$$
\begin{aligned}
\left|T_{f i}\right|^{2}= & \frac{1}{2}\left[\operatorname{Tr}\left(\gamma_{\mu} \not p_{i} \gamma_{\nu} \not p_{f}\right) J_{\mu} J_{\nu}^{*}+m^{2} \gamma_{\mu} \gamma_{\nu} J_{\mu} J_{\nu}^{*}\right] \\
= & \frac{1}{2}\left[\operatorname{Tr}\left(\gamma_{\mu} \gamma_{\rho} \gamma_{\nu} \gamma_{\lambda}\right)\left(p_{i}\right)_{\rho}\left(p_{f}\right)_{\lambda} J_{\mu} J_{\nu}^{*}+m^{2} \operatorname{Tr}\left(\gamma_{\mu} \gamma_{\nu}\right) J_{\mu} J_{\nu}^{*}\right] \\
= & \frac{1}{2}\left[\left\{4 g_{\mu \lambda} g_{\nu \rho}-4 g_{\mu \nu} g_{\lambda \rho}+4 g_{\mu \rho} g_{\nu \lambda}\right\}\left(p_{i}\right)_{\rho}\left(p_{f}\right)_{\lambda} J_{\mu} J_{\nu}^{*}\right. \\
& \left.+4 m^{2}\left(g_{\mu \nu}\right) J_{\mu} J_{\nu}^{*}\right] \\
= & \frac{1}{2}\left[4\left(\mathbf{p}_{f} \cdot \mathbf{J}\right)\left(\mathbf{p}_{i} \cdot \mathbf{J}^{*}\right)-4\left(\mathbf{p}_{i} \cdot \mathbf{p}_{f}\right)\left(\mathbf{J} \cdot \mathbf{J}^{*}\right)\right. \\
& \left.+4\left(\mathbf{p}_{i} \cdot \mathbf{J}\right)\left(\mathbf{p}_{f} \cdot \mathbf{J}^{*}\right)+4 m^{2}\left(\mathbf{J} \cdot \mathbf{J}^{*}\right)\right] \\
= & 2\left[\left(\mathbf{p}_{i} \cdot \mathbf{J}\right)\left(\mathbf{p}_{f} \cdot \mathbf{J}^{*}\right)+\left(\mathbf{p}_{f} \cdot \mathbf{J}\right)\left(\mathbf{p}_{i} \cdot \mathbf{J}^{*}\right)+\left(m^{2}-\mathbf{p}_{i} \cdot \mathbf{p}_{f}\right) \mathbf{J} \cdot \mathbf{J}^{*}\right]
\end{aligned}
$$

where $\left(\mathbf{p}_{i} \cdot \mathbf{J}\right), \ldots$ are scalar products of four-vectors.
2.5 Given

$$
\mathcal{O}=\phi \not b \phi
$$

then

$$
\begin{aligned}
\tilde{\mathcal{O}} & =\gamma_{0} \mathcal{O}^{\dagger} \gamma_{0} \\
& =\gamma_{0}(d b b)^{\dagger} \gamma_{0} \\
& =\gamma_{0}\left(\gamma_{\mu} a_{\mu} \gamma_{\nu} b_{\nu} \gamma_{\lambda} c_{\lambda}\right)^{\dagger} \gamma_{0} \\
& =\gamma_{0}\left\{\left(\gamma_{\mu} \gamma_{\nu} \gamma_{\lambda}\right)^{\dagger} a_{\mu} b_{\nu} c_{\lambda}\right\} \gamma_{0} \\
& =\gamma_{0}\left(\gamma_{\lambda}^{\dagger} \gamma_{\nu}^{\dagger} \gamma_{\mu}^{\dagger} a_{\mu} b_{\nu} c_{\lambda}\right) \gamma_{0} \\
& =\gamma_{0}\left\{\left(\gamma_{0} c_{0}+\boldsymbol{\gamma} \cdot \boldsymbol{c}\right)\left(\gamma_{0} b_{0}+\boldsymbol{\gamma} \cdot \boldsymbol{b}\right)\left(\gamma_{0} a_{0}+\boldsymbol{\gamma} \cdot \boldsymbol{a}\right)\right\} \gamma_{0}
\end{aligned}
$$

since $\gamma_{0}^{\dagger}=\gamma_{0}$ and $\gamma_{k}^{\dagger}=-\gamma_{k}, k=1,2,3$. Switching the $\gamma_{0}$ in the front to the rear side using the commutation relations $\gamma_{0} \gamma=-\gamma \gamma_{0}$ and noting that $\gamma_{0}^{2}=1$, we get

$$
\begin{aligned}
\tilde{\mathcal{O}} & =\left(\gamma_{0} c_{0}-\boldsymbol{\gamma} \cdot \boldsymbol{c}\right)\left(\gamma_{0} b_{0}-\boldsymbol{\gamma} \cdot \boldsymbol{b}\right)\left(\gamma_{0} a_{0}-\boldsymbol{\gamma} \cdot \boldsymbol{a}\right) \\
& =\dot{b} b \boldsymbol{d}
\end{aligned}
$$

## Chapter 3

## The Dirac Equation with External Potentials

Having considered the free particle Dirac equation and its solutions in the last chapter, let us consider now the Dirac equation with external potentials. Let us investigate how the charged spin- $\frac{1}{2}$ particle interacts with the external fields such as electro-magnetic field. Even before the advent of Dirac's theory, this was studied extensively using the Schrödinger equation along with the concept of spin, as introduced by Pauli, in order to explain the experimental observations. In Dirac's theory, the spin comes out as a natural consequence and let us examine whether an effective Hamiltonian can be obtained in the non-relativistic limit, that can be compared with the earlier theory that was developed with ad hoc inputs. The effective Hamiltonian obtained from Dirac's equation includes a magnetic moment term with the correct gyro-magnetic ratio for the electron spin, spin-orbit coupling term and the Darwin term. It is also shown how the Dirac equation can be solved for a spherically symmetric potential such as the Coulomb potential to obtain the energy levels of hydrogen-like atom with the fine structure splitting arising from spin-orbit interaction in agreement with the experimental observation.

### 3.1 Dirac Equation with Electromagnetic Field

The Dirac equation for a spin- $\frac{1}{2}$ particle of charge $e$ in an electro-magnetic field, characterized by the vector and scalar potentials $\boldsymbol{A}$ and $\phi$, is given
by ${ }^{1}$

$$
\begin{equation*}
\left\{c \boldsymbol{\alpha} \cdot\left(\boldsymbol{p}-\frac{e \boldsymbol{A}}{c}\right)+\beta m c^{2}\right\} \psi=(E-e \phi) \psi \tag{3.1}
\end{equation*}
$$

The electric and magnetic field strengths of the electro-magnetic field are defined in terms of the vector and scalar potentials.

$$
\begin{equation*}
\boldsymbol{E}=-\frac{1}{c} \frac{\partial \boldsymbol{A}}{\partial t}-\nabla \phi ; \quad \boldsymbol{B}=\operatorname{curl} \boldsymbol{A}=\boldsymbol{\nabla} \times \boldsymbol{A} . \tag{3.2}
\end{equation*}
$$

### 3.2 Non-relativistic approximation to Dirac equation

### 3.2.1 Constant magnetic field

Let $\phi=0$ and $\boldsymbol{P}=\boldsymbol{p}-\frac{e \boldsymbol{A}}{c}$. Then Eq. (3.1) becomes

$$
\begin{equation*}
\left(c \boldsymbol{\alpha} \cdot \boldsymbol{P}+\beta m c^{2}\right) \psi=E \psi \tag{3.3}
\end{equation*}
$$

The state vector $\psi$ is a four component spinor but we shall use a shorthand notation as two-component spinor. Now, Eq. (3.3) can be written as a $2 \times 2$ matrix equation

$$
\left\{c\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \boldsymbol{P}  \tag{3.4}\\
\boldsymbol{\sigma} \cdot \boldsymbol{P} & 0
\end{array}\right]+m c^{2}\left[\begin{array}{cc}
I & 0 \\
0 & -I
\end{array}\right]\right\}\left[\begin{array}{c}
\varphi_{a} \\
\varphi_{b}
\end{array}\right]=E\left[\begin{array}{c}
\varphi_{a} \\
\varphi_{b}
\end{array}\right] .
$$

The matrix equation (3.4) yields two coupled equations in $\varphi_{a}$ and $\varphi_{b}$.

$$
\begin{align*}
c \boldsymbol{\sigma} \cdot \boldsymbol{P} \varphi_{b}+\left(m c^{2}-E\right) \varphi_{a} & =0 ;  \tag{3.5}\\
c \boldsymbol{\sigma} \cdot \boldsymbol{P} \varphi_{a}-\left(m c^{2}+E\right) \varphi_{b} & =0 . \tag{3.6}
\end{align*}
$$

We are interested in the positive energy states with $\varphi_{a}$ as large component and $\varphi_{b}$ as small component. If $\epsilon$ is the kinetic energy of the electron, then

$$
E=m c^{2}+\epsilon,
$$

[^14]where $m$ is the rest mass of the charged particle. For small $\epsilon$, one can go to the non-relativistic limit and investigate the Dirac equation in the presence of constant magnetic field. From Eq. (3.6), we get
\[

$$
\begin{equation*}
\varphi_{b}=\frac{c \boldsymbol{\sigma} \cdot \boldsymbol{P}}{m c^{2}+E} \varphi_{a} \approx \frac{c \boldsymbol{\sigma} \cdot \boldsymbol{P}}{2 m c^{2}} \varphi_{a} \tag{3.7}
\end{equation*}
$$

\]

$\varphi_{b}$ is of the order of $(v / c) \varphi_{a}$ and hence it is referred to as a small component. Substituting (3.7) in Eq. (3.5), we get

$$
\begin{equation*}
\left\{\frac{(\boldsymbol{\sigma} \cdot \boldsymbol{P})^{2}}{2 m}+\left(m c^{2}-E\right)\right\} \varphi_{a}=0 \quad \text { or } \quad \frac{(\boldsymbol{\sigma} \cdot \boldsymbol{P})^{2}}{2 m} \varphi_{a}=\epsilon \varphi_{a} \tag{3.8}
\end{equation*}
$$

The above equation is the non-relativistic limit of Dirac equation and it is an equation in $\varphi_{a}$. Using the vector algebra, we get

$$
\begin{align*}
(\boldsymbol{\sigma} \cdot \boldsymbol{P})^{2} & =\boldsymbol{P} \cdot \boldsymbol{P}+i \boldsymbol{\sigma} \cdot(\boldsymbol{P} \times \boldsymbol{P}),  \tag{3.9}\\
\boldsymbol{P} \times \boldsymbol{P} & =\left(\boldsymbol{p}-\frac{e \boldsymbol{A}}{c}\right) \times\left(\boldsymbol{p}-\frac{e \boldsymbol{A}}{c}\right) \\
& =-\frac{e}{c}(\boldsymbol{A} \times \boldsymbol{p}+\boldsymbol{p} \times \boldsymbol{A})=i \hbar \frac{e}{c} \boldsymbol{\nabla} \times \boldsymbol{A} . \tag{3.10}
\end{align*}
$$

Substituting (3.9) and (3.10) into Eq. (3.8) and observing that the strength of the magnetic field $\boldsymbol{B}=\boldsymbol{\nabla} \times \boldsymbol{A}$, we get

$$
\begin{equation*}
\left\{\frac{1}{2 m}\left(\boldsymbol{p}-\frac{e \boldsymbol{A}}{c}\right)^{2}-\frac{e \hbar}{2 m c} \boldsymbol{\sigma} \cdot \boldsymbol{B}\right\} \varphi_{a}=\epsilon \varphi_{a} . \tag{3.11}
\end{equation*}
$$

One can recognize that the second term on the left hand side of the above equation as arising from the magnetic dipole moment

$$
\begin{equation*}
\boldsymbol{M}=\frac{e \hbar}{2 m c} \boldsymbol{\sigma}=\mu_{B} \boldsymbol{\sigma} \tag{3.12}
\end{equation*}
$$

where $\mu_{B}=e \hbar / 2 m c$ is known as the Bohr magneton.
It is possible to expand the first term on the left hand side of Eq. (3.11) to get more physical insight.

$$
\begin{align*}
\left(\boldsymbol{p}-\frac{e \boldsymbol{A}}{c}\right)^{2} & =p^{2}+\frac{e^{2} A^{2}}{c^{2}}-\frac{e}{c}(\boldsymbol{p} \cdot \boldsymbol{A}+\boldsymbol{A} \cdot \boldsymbol{p}) \\
& =p^{2}+\frac{e^{2} A^{2}}{c^{2}}+\frac{i e \hbar}{c} \boldsymbol{\nabla} \cdot \boldsymbol{A}-\frac{2 e}{c} \boldsymbol{A} \cdot \boldsymbol{p} . \tag{3.13}
\end{align*}
$$

The last step in Eq. (3.13) is obtained by treating $\boldsymbol{p}$ as a differential operator $-i \hbar \nabla$.

$$
\boldsymbol{p} \cdot \boldsymbol{A}=-i \hbar(\boldsymbol{\nabla} \cdot \boldsymbol{A}+\boldsymbol{A} \cdot \boldsymbol{\nabla})=-i \hbar \boldsymbol{\nabla} \cdot \boldsymbol{A}+\boldsymbol{A} \cdot \boldsymbol{p}
$$

Further, for constant magnetic field ${ }^{2}$,

$$
\boldsymbol{A}=\frac{1}{2} \boldsymbol{B} \times \boldsymbol{r}
$$

such that

$$
\begin{equation*}
\boldsymbol{A} \cdot \boldsymbol{p}=\frac{1}{2} \boldsymbol{B} \times \boldsymbol{r} \cdot \boldsymbol{p}=\frac{1}{2} \boldsymbol{B} \cdot(\boldsymbol{r} \times \boldsymbol{p})=\frac{1}{2} \boldsymbol{B} \cdot \boldsymbol{L} . \tag{3.14}
\end{equation*}
$$

Substituting the results (3.13) and (3.14) into Eq. (3.11), we get

$$
\left\{\frac{p^{2}}{2 m}+\frac{e^{2} A^{2}}{2 m c^{2}}+\frac{i e \hbar}{2 m c} \boldsymbol{\nabla} \cdot \boldsymbol{A}-\frac{e}{m c} \boldsymbol{B} \cdot(\boldsymbol{L}+\hbar \boldsymbol{\sigma})\right\} \varphi_{a}=\epsilon \varphi_{a} .
$$

Since the spin angular momentum $\boldsymbol{S}=\frac{1}{2} \hbar \boldsymbol{\sigma}$, we can rewrite the above equation as

$$
\begin{equation*}
\left\{\frac{p^{2}}{2 m}+\frac{e^{2} A^{2}}{2 m c^{2}}+\frac{i e \hbar}{2 m c} \boldsymbol{\nabla} \cdot \boldsymbol{A}-\frac{e}{m c} \boldsymbol{B} \cdot(\boldsymbol{L}+2 \boldsymbol{S})\right\} \varphi_{a}=\epsilon \varphi_{a} . \tag{3.15}
\end{equation*}
$$

The last term on the left hand side of Eq. (3.15) is due to the interaction of orbital and spin angular momentum with the external magnetic field $\boldsymbol{B}$. It is remarkable that the Dirac equation yields the gyromagnetic ratio for spin as 2 .

### 3.2.2 Static electric field or Electrostatic field

Let us now consider a spin- $\frac{1}{2}$ particle of charge $e$ in an electrostatic field of strength $\mathcal{E}$. In this case, the vector potential $\boldsymbol{A}=0$.

$$
\begin{equation*}
e \mathcal{E}=-\boldsymbol{\nabla} V, \quad V=e \phi . \tag{3.16}
\end{equation*}
$$

[^15]$$
\boldsymbol{A}=\frac{1}{2} \boldsymbol{B} \times \boldsymbol{r}, \quad\left(A_{x}=\frac{1}{2}\left(B_{y} z-B_{z} y\right), A_{y}=\frac{1}{2}\left(B_{z} x-B_{x} z\right), A_{z}=\frac{1}{2}\left(B_{x} y-B_{y} x\right)\right)
$$
satisfies the Eq. curl $\boldsymbol{A}=\boldsymbol{B}$.
$$
(\operatorname{curl} \boldsymbol{A})_{x}=(\boldsymbol{\nabla} \times \boldsymbol{A})_{x}=\frac{\partial A_{z}}{\partial y}-\frac{\partial A_{y}}{\partial z}=B_{x} .
$$

The Dirac equation for the electrostatic potential can be written as

$$
\begin{equation*}
\left(c \boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m c^{2}\right) \psi=(E-e \phi) \psi . \tag{3.17}
\end{equation*}
$$

This can be explicitly written in the $2 \times 2$ matrix form.

$$
\left\{c\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \boldsymbol{p}  \tag{3.18}\\
\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0
\end{array}\right]+m c^{2}\left[\begin{array}{cc}
I & 0 \\
0 & -I
\end{array}\right]\right\}\left[\begin{array}{l}
\varphi_{a} \\
\varphi_{b}
\end{array}\right]=(E-V)\left[\begin{array}{l}
\varphi_{a} \\
\varphi_{b}
\end{array}\right] .
$$

This leads to two coupled equations in $\varphi_{a}$ and $\varphi_{b}$.

$$
\begin{align*}
\left(E-m c^{2}-V\right) \varphi_{a} & =c \boldsymbol{\sigma} \cdot \boldsymbol{p} \varphi_{b}  \tag{3.19}\\
\left(E+m c^{2}-V\right) \varphi_{b} & =c \boldsymbol{\sigma} \cdot \boldsymbol{p} \varphi_{a} \tag{3.20}
\end{align*}
$$

From Eq. (3.20), we obtain

$$
\begin{equation*}
\varphi_{b}=\frac{c \boldsymbol{\sigma} \cdot \boldsymbol{p}}{E+m c^{2}-V} \varphi_{a}=\frac{c \boldsymbol{\sigma} \cdot \boldsymbol{p}}{2 m c^{2}+\epsilon-V} \varphi_{a} \tag{3.21}
\end{equation*}
$$

with $E=\epsilon+m c^{2}$. The total energy $E$ is written as a sum of kinetic energy $(\epsilon)$ and rest mass energy $\left(m c^{2}\right)$. Substituting (3.21) into (3.19), we obtain an equation for $\varphi_{a}$.

$$
\begin{align*}
(\epsilon-V) \varphi_{a} & =c(\boldsymbol{\sigma} \cdot \boldsymbol{p}) \frac{c}{2 m c^{2}+\epsilon-V}(\boldsymbol{\sigma} \cdot \boldsymbol{p}) \varphi_{a} \\
& =c(\boldsymbol{\sigma} \cdot \boldsymbol{p}) \frac{c}{2 m c^{2}\left(1+\frac{\epsilon-V}{2 m c^{2}}\right)}(\boldsymbol{\sigma} \cdot \boldsymbol{p}) \varphi_{a} \\
& =\frac{(\boldsymbol{\sigma} \cdot \boldsymbol{p})}{2 m}\left(1-\frac{\epsilon-V}{2 m c^{2}}\right)(\boldsymbol{\sigma} \cdot \boldsymbol{p}) \varphi_{a} \\
& =\left\{\frac{p^{2}}{2 m}-\frac{\epsilon p^{2}}{4 m^{2} c^{2}}+\frac{1}{4 m^{2} c^{2}}(\boldsymbol{\sigma} \cdot \boldsymbol{p}) V(\boldsymbol{\sigma} \cdot \boldsymbol{p})\right\} \varphi_{a} . \tag{3.22}
\end{align*}
$$

In Dirac equation, $\varphi_{a}$ is only the large component of $\psi$ and so only $\psi$ is normalized. So, a correction has to be made for normalization, when we are using only $\varphi_{a}$.

$$
\begin{align*}
\int \psi^{*} \psi d^{3} x & =\int\left(\left|\varphi_{a}\right|^{2}+\left|\varphi_{b}\right|^{2}\right) d^{3} x \\
& =\int\left\{1+\left(\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{2 m c}\right)^{2}\right\}\left|\varphi_{a}\right|^{2} d^{3} x \\
& =\left\{1+\frac{p^{2}}{4 m^{2} c^{2}}\right\}\left|\varphi_{a}\right|^{2} d^{3} x \\
& =\left|\varphi_{N}\right|^{2} d^{3} x=1 \tag{3.23}
\end{align*}
$$

where

$$
\begin{equation*}
\left|\varphi_{N}\right|^{2}=\left(1+\frac{p^{2}}{4 m^{2} c^{2}}\right)\left|\varphi_{a}\right|^{2} . \tag{3.24}
\end{equation*}
$$

In deducing Eq. (3.23), we have used an approximate expression for $\varphi_{b} \approx \frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{2 m c} \varphi_{a}$ deduced from Eq. (3.21), neglecting the terms $\epsilon$ and $V$ which are very small compared to $2 m c^{2}$. From Eq. (3.24), we get

$$
\begin{equation*}
\varphi_{a}=\frac{\varphi_{N}}{\left(1+\frac{p^{2}}{4 m^{2} c^{2}}\right)^{1 / 2}}=\left(1-\frac{p^{2}}{8 m^{2} c^{2}}\right) \varphi_{N} \tag{3.25}
\end{equation*}
$$

Substituting (3.25) in Eq. (3.22), we get

$$
\begin{align*}
\{ & \left.\epsilon-V-\frac{\epsilon p^{2}}{8 m^{2} c^{2}}+\frac{V p^{2}}{8 m^{2} c^{2}}\right\} \varphi_{N} \\
& =\left\{\frac{p^{2}}{2 m}-\frac{\epsilon p^{2}}{4 m^{2} c^{2}}+\frac{1}{4 m^{2} c^{2}}(\boldsymbol{\sigma} \cdot \boldsymbol{p}) V(\boldsymbol{\sigma} \cdot \boldsymbol{p})\right\}\left(1-\frac{p^{2}}{8 m^{2} c^{2}}\right) \varphi_{N} \\
& =\left\{\frac{p^{2}}{2 m}-\frac{p^{4}}{16 m^{3} c^{2}}-\frac{\epsilon p^{2}}{4 m^{2} c^{2}}+\frac{1}{4 m^{2} c^{2}}(\boldsymbol{\sigma} \cdot \boldsymbol{p}) V(\boldsymbol{\sigma} \cdot \boldsymbol{p})\right\} \varphi_{N}, \tag{3.26}
\end{align*}
$$

retaining only terms up to the order $\frac{1}{m^{3} c^{2}}$. Simplifying Eq. (3.26), we get

$$
\begin{align*}
\epsilon \varphi_{N}= & \left\{\frac{p^{2}}{2 m}+V-\frac{p^{4}}{16 m^{3} c^{2}}-\frac{\epsilon p^{2}}{8 m^{2} c^{2}}-\frac{V p^{2}}{8 m^{2} c^{2}}\right. \\
& \left.+\frac{1}{4 m^{2} c^{2}}(\boldsymbol{\sigma} \cdot \boldsymbol{p}) V(\boldsymbol{\sigma} \cdot \boldsymbol{p})\right\} \varphi_{N}, \tag{3.27}
\end{align*}
$$

In Eq. (3.27), the non-relativistic energy $\epsilon$ occurs also on the right-hand side; for which one can substitute an approximate expression.

$$
\begin{equation*}
\epsilon p^{2}=p^{2} \epsilon=p^{2}\left(\frac{p^{2}}{2 m}+V\right)=\frac{p^{4}}{2 m}+p^{2} V . \tag{3.28}
\end{equation*}
$$

Substituting (3.28) into Eq. (3.27) and simplifying, we can get an expression for the effective non-relativistic Hamiltonian $H_{\text {eff }}$. Denoting $\epsilon$ on the left hand side by $H_{\text {eff }}$, we get

$$
\begin{align*}
H_{\mathrm{eff}}= & \frac{p^{2}}{2 m}+V-\frac{p^{4}}{8 m^{3} c^{2}}-\frac{1}{8 m^{2} c^{2}}\left(p^{2} V+V p^{2}\right) \\
& +\frac{1}{4 m^{2} c^{2}}(\boldsymbol{\sigma} \cdot \boldsymbol{p}) V(\boldsymbol{\sigma} \cdot \boldsymbol{p}) \tag{3.29}
\end{align*}
$$

Remembering that $\boldsymbol{p}$ is an operator $-i \hbar \boldsymbol{\nabla}$, we find

$$
\begin{align*}
(\boldsymbol{\sigma} \cdot \boldsymbol{p}) V(\boldsymbol{\sigma} \cdot \boldsymbol{p}) & =-i \hbar(\boldsymbol{\sigma} \cdot \boldsymbol{\nabla}) V(\boldsymbol{\sigma} \cdot \boldsymbol{p}) \\
& =-i \hbar(\boldsymbol{\sigma} \cdot \boldsymbol{\nabla} V)(\boldsymbol{\sigma} \cdot \boldsymbol{p})+V(\boldsymbol{\sigma} \cdot \boldsymbol{p})(\boldsymbol{\sigma} \cdot \boldsymbol{p}) \\
& =-i \hbar\{(\boldsymbol{\nabla} V \cdot \boldsymbol{p})+i \boldsymbol{\sigma} \cdot \boldsymbol{\nabla} V \times \boldsymbol{p}\}+V p^{2} \tag{3.30}
\end{align*}
$$

If $\mathcal{E}$ is the electric field in which the particle is moving, then

$$
\begin{equation*}
e \mathcal{E}=-\nabla V=-\frac{1}{r} \frac{d V}{d r} \boldsymbol{r} \tag{3.31}
\end{equation*}
$$

Using (3.31), we get

$$
\begin{align*}
(\boldsymbol{\sigma} \cdot \boldsymbol{p}) V(\boldsymbol{\sigma} \cdot \boldsymbol{p}) & =i \hbar e(\mathcal{E} \cdot \boldsymbol{p})-\hbar e \boldsymbol{\sigma} \cdot \boldsymbol{\mathcal { E }} \times \boldsymbol{p}+V p^{2} \\
& =i \hbar e(\mathcal{E} \cdot \boldsymbol{p})+\hbar e \frac{1}{r} \frac{d V}{d r} \boldsymbol{\sigma} \cdot \boldsymbol{r} \times \boldsymbol{p}+V p^{2} \\
& =i \hbar e(\mathcal{E} \cdot \boldsymbol{p})+\hbar e \frac{1}{r} \frac{d V}{d r} \boldsymbol{\sigma} \cdot \boldsymbol{L}+V p^{2} \\
& =i \hbar e(\mathcal{E} \cdot \boldsymbol{p})+e \frac{2}{r} \frac{d V}{d r} \boldsymbol{S} \cdot \boldsymbol{L}+V p^{2} . \tag{3.32}
\end{align*}
$$

Substituting (3.32) in Eq. (3.29) and simplifying, we get

$$
\begin{align*}
H_{\mathrm{eff}}= & \frac{p^{2}}{2 m}+V-\frac{p^{4}}{8 m^{3} c^{2}}-\frac{1}{8 m^{2} c^{2}}\left(p^{2} V-V p^{2}\right) \\
& +\frac{1}{4 m^{2} c^{2}}\left\{i \hbar e(\boldsymbol{\mathcal { E }} \cdot \boldsymbol{p})+\frac{2 e}{r} \frac{d V}{d r}(\boldsymbol{S} \cdot \boldsymbol{L})\right\} \tag{3.33}
\end{align*}
$$

Since $\boldsymbol{p}$ is an operator, we get

$$
\boldsymbol{p} V=-i \hbar(\boldsymbol{\nabla} V+V \boldsymbol{\nabla})=i \hbar e \mathcal{E}+V \boldsymbol{p}
$$

Equivalently,

$$
[\boldsymbol{p}, V]_{-}=i \hbar e \mathcal{E}
$$

It follows that

$$
\begin{align*}
{\left[\boldsymbol{p}^{2}, V\right]_{-} } & =\boldsymbol{p}[\boldsymbol{p}, V]_{-}+[\boldsymbol{p}, V]_{-} \boldsymbol{p} \\
& =i \hbar e(\boldsymbol{p} \cdot \mathcal{E}+\mathcal{E} \cdot \boldsymbol{p}) \tag{3.34}
\end{align*}
$$

Substituting (3.34) in Eq. (3.33), we get

$$
\begin{align*}
H_{\mathrm{eff}}= & \frac{p^{2}}{2 m}+V-\frac{p^{4}}{8 m^{3} c^{2}}-\frac{i \hbar e}{8 m^{2} c^{2}}(\boldsymbol{p} \cdot \boldsymbol{\mathcal { E }}+\boldsymbol{\mathcal { E }} \cdot \boldsymbol{p}) \\
& +\frac{1}{4 m^{2} c^{2}}\left\{i \hbar e(\mathcal{E} \cdot \boldsymbol{p})+\frac{2 e}{r} \frac{d V}{d r}(\boldsymbol{S} \cdot \boldsymbol{L})\right\} \\
= & \frac{p^{2}}{2 m}+V-\frac{p^{4}}{8 m^{3} c^{2}}-\frac{i \hbar e}{8 m^{2} c^{2}}(\boldsymbol{p} \cdot \boldsymbol{\mathcal { E }}-\boldsymbol{\mathcal { E }} \cdot \boldsymbol{p})+\frac{1}{4 m^{2} c^{2}} \frac{2 e}{r} \frac{d V}{d r}(\boldsymbol{S} \cdot \boldsymbol{L}) \\
= & \frac{p^{2}}{2 m}+V-\frac{p^{4}}{8 m^{3} c^{2}}-\frac{\hbar^{2} e}{8 m^{2} c^{2}} \boldsymbol{\nabla} \cdot \mathcal{E}+\frac{1}{4 m^{2} c^{2}} \frac{2 e}{r} \frac{d V}{d r}(\boldsymbol{S} \cdot \boldsymbol{L}) . \tag{3.35}
\end{align*}
$$

The last step is obtained by treating $\boldsymbol{p}$ as a gradient operator.

$$
\boldsymbol{p} \cdot \mathcal{E}=-i \hbar \boldsymbol{\nabla} \cdot \mathcal{E}+\mathcal{E} \cdot \boldsymbol{p}, \quad \text { such that } \quad \boldsymbol{p} \cdot \mathcal{E}-\mathcal{E} \cdot \boldsymbol{p}=-i \hbar \boldsymbol{\nabla} \cdot \mathcal{E}
$$

Let us examine the effective Hamiltonian (3.35) that we have obtained as the non-relativistic limit of Dirac equation in static electric potential. The first term is the kinetic energy term, the second term is the scalar potential energy term, the third term can be considered as the relativistic correction ${ }^{3}$ to the kinetic energy, the fourth term is the Darwin term and the fifth term is the spin-orbit coupling term.

Historically, the effective Hamiltonian (3.35) was obtained much before the discovery of Dirac equation and it is the greatest triumph of Dirac equation that it yields the same effective Hamiltonian in the nonrelativistic limit.

### 3.3 The Dirac equation for a central potential

Let us consider the Dirac equation for a central potential (in natural units with $\hbar=c=1$ ).

$$
\begin{equation*}
\{\boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m+V(r)\} \psi(\boldsymbol{r})=E \psi(\boldsymbol{r}), \tag{3.36}
\end{equation*}
$$

where the potential $V(r)$ is spherically symmetric and depends only on the radial coordinate $r$ and not on its angular coordinates $\theta, \phi$.

[^16]The Dirac Hamiltonian $H$

$$
\begin{equation*}
H=\boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m+V(r) \tag{3.37}
\end{equation*}
$$

does not commute with the orbital angular momentum operator $\boldsymbol{L}$ but commutes with the total angular momentum operator $\boldsymbol{J}=\boldsymbol{L}+\boldsymbol{S}$ and another operator $K$ defined by

$$
\begin{equation*}
K=\beta(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1) \tag{3.38}
\end{equation*}
$$

Also, the operators $\boldsymbol{J}$ and $K$ commute among themselves (vide Solved Problem 2.2).

$$
\begin{equation*}
[H, \boldsymbol{J}]_{-}=0 ; \quad\left[H, J_{z}\right]_{-}=0 ; \quad[H, K]_{-}=0 ; \quad[\boldsymbol{J}, K]_{-}=0 \tag{3.39}
\end{equation*}
$$

So, it is possible to find the simultaneous eigenfunctions and eigenvalues of operators $H, \boldsymbol{J}, J_{z}$ and $K$. Thus the solutions for the Dirac equation with a central potential are characterized by the quantum numbers $j, m, \kappa$.

$$
\begin{align*}
H \psi_{\kappa, j, m}(\boldsymbol{r}) & =E_{\kappa, j, m} \psi_{\kappa, j, m}(\boldsymbol{r}) ;  \tag{3.40}\\
\boldsymbol{J}^{2} \psi_{\kappa, j, m}(\boldsymbol{r}) & =j(j+1) \psi_{\kappa, j . m}(\boldsymbol{r}) ;  \tag{3.41}\\
J_{z} \psi_{\kappa, j, m}(\boldsymbol{r}) & =m \psi_{\kappa, j . m}(\boldsymbol{r}) ;  \tag{3.42}\\
K \psi_{\kappa, j, m}(\boldsymbol{r}) & =\kappa \psi_{\kappa, j \cdot m}(\boldsymbol{r}) . \tag{3.43}
\end{align*}
$$

## Eigenvalues of the operator $K$

Squaring the operator $K$, we get

$$
\begin{align*}
K^{2} & =\beta(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1) \beta(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1) \\
& =(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1)(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1) \\
& =1+2 \boldsymbol{\sigma} \cdot \boldsymbol{L}+(\boldsymbol{\sigma} \cdot \boldsymbol{L})(\boldsymbol{\sigma} \cdot \boldsymbol{L}) \\
& =1+\boldsymbol{\sigma} \cdot \boldsymbol{L}+\boldsymbol{L}^{2} \tag{3.44}
\end{align*}
$$

since $\beta$ commutes with $\boldsymbol{\sigma} \cdot \boldsymbol{L}$ and $\beta^{2}=1$. Besides,

$$
(\boldsymbol{\sigma} \cdot \boldsymbol{L})^{2}=\boldsymbol{L}^{2}+i \boldsymbol{\sigma} \cdot(\boldsymbol{L} \times \boldsymbol{L})=\boldsymbol{L}^{2}-\boldsymbol{\sigma} \cdot \boldsymbol{L}, \quad \text { since } \boldsymbol{L} \times \boldsymbol{L}=i \boldsymbol{L} .
$$

Thus

$$
\begin{equation*}
K^{2}=\left(\boldsymbol{L}+\frac{1}{2} \boldsymbol{\sigma}\right)^{2}-\frac{1}{4} \boldsymbol{\sigma}^{2}+1=\boldsymbol{J}^{2}+\frac{1}{4}, \tag{3.45}
\end{equation*}
$$

since $\boldsymbol{J}=\boldsymbol{L}+\frac{1}{2} \boldsymbol{\sigma}$ and $\boldsymbol{\sigma}^{2}=3$.

Thus, we find that a Dirac particle in central potential such as Hydrogen atom has eigenstates which are not only the eigenfunctions of the operator $\boldsymbol{J}^{2}$ with eigenvalues $j(j+1)$ but also the eigenfunctions of the operator $K$ with eigenvalue $\kappa$.

$$
\begin{equation*}
\kappa^{2}=j(j+1)+\frac{1}{4}=\left(j+\frac{1}{2}\right)^{2} \quad \text { or } \quad \kappa= \pm\left(j+\frac{1}{2}\right) . \tag{3.46}
\end{equation*}
$$

Since $j$ can take positive half-odd integral values $\left(j=\frac{1}{2}, \frac{3}{2}, \frac{5}{2}, \cdots\right)$, the eigenvalues of $K$ are $\kappa= \pm 1, \pm 2, \pm 3, \cdots$, excluding zero.

The Hamiltonian $H$ is a $4 \times 4$ matrix $^{4}$ and so the eigenfunction is a four-component column vector. But it is found more convenient to write the Hamiltonian in a $2 \times 2$ matrix form, treating each element of which as a $2 \times 2$ matrix. In a similar way, the eigenfunction $\psi(\boldsymbol{r})$ can be written in the two component form

$$
\psi(\boldsymbol{r})=\left[\begin{array}{c}
\Phi_{a}(\boldsymbol{r})  \tag{3.47}\\
\Phi_{b}(\boldsymbol{r})
\end{array}\right]
$$

treating each element, in turn, as a two-component column vector known as spinor.

### 3.3.1 Dirac Hamiltonian in spherical polar coordinates

Of the three terms in the Hamiltonian (3.37), only the first term requires consideration, since the other two terms do not have any angular dependence.

$$
\begin{equation*}
\boldsymbol{\alpha} \cdot \boldsymbol{p}=-i \boldsymbol{\alpha} \cdot \nabla \tag{3.48}
\end{equation*}
$$

A convenient form for the gradient operator $\boldsymbol{\nabla}$ can be obtained from the expansion of a triple vector product.

$$
\begin{equation*}
a \times(b \times c)=b(a \cdot c)-c(a \cdot b) \tag{3.49}
\end{equation*}
$$

Substituting $\boldsymbol{a}=\boldsymbol{b}=\hat{\boldsymbol{r}}$ and $\boldsymbol{c}=\boldsymbol{\nabla}$, we get on rearrangement

$$
\begin{align*}
\boldsymbol{\nabla} & =\hat{\boldsymbol{r}}(\hat{\boldsymbol{r}} \cdot \boldsymbol{\nabla})-\hat{\boldsymbol{r}} \times(\hat{\boldsymbol{r}} \times \boldsymbol{\nabla}) \\
& =\hat{\boldsymbol{r}} \frac{\partial}{\partial r}-\frac{i}{r}\{\hat{\boldsymbol{r}} \times(\boldsymbol{r} \times \boldsymbol{p})\} \\
& =\hat{\boldsymbol{r}} \frac{\partial}{\partial r}-\frac{i}{r}(\hat{\boldsymbol{r}} \times \boldsymbol{L}), \tag{3.50}
\end{align*}
$$

[^17]where $\hat{\boldsymbol{r}} \equiv \hat{\boldsymbol{e}}_{r}$ is the unit radius vector and $\boldsymbol{L}$, the orbital angular momentum operator. Substituting (3.50) into Eq. (3.48), we get
\[

$$
\begin{align*}
\boldsymbol{\alpha} \cdot \boldsymbol{p} & =-i(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}}) \frac{\partial}{\partial r}-\frac{1}{r} \boldsymbol{\alpha} \cdot(\hat{\boldsymbol{r}} \times \boldsymbol{L}) \\
& =-i(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}}) \frac{\partial}{\partial r}+\frac{i}{r}(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \boldsymbol{L}) \tag{3.51}
\end{align*}
$$
\]

using the relation (vide Solved Problem 2.3)

$$
\begin{equation*}
\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}} \times \boldsymbol{L}=-i(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \boldsymbol{L}) . \tag{3.52}
\end{equation*}
$$

Using the definition (3.38) of the operator $K$, we can express the operator $\boldsymbol{\sigma} \cdot \boldsymbol{L}$ in terms of $K$, by multiplying both sides of Eq. (3.38) by $\beta$ and using the property $\beta^{2}=1$.

$$
\begin{equation*}
\boldsymbol{\sigma} \cdot \boldsymbol{L}=\beta K-1 \tag{3.53}
\end{equation*}
$$

Substituting (3.53) into Eq. (3.51), we get

$$
\begin{equation*}
\boldsymbol{\alpha} \cdot \boldsymbol{p}=-i(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}}) \frac{\partial}{\partial r}+\frac{i}{r}(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})(\beta K-1) . \tag{3.54}
\end{equation*}
$$

Using (3.54), we can rewrite the Dirac equation (3.36) for spherically symmetric potential in a form, suitable for solution in spherical polar coordinates ${ }^{5}$.

$$
\begin{equation*}
\left[-i(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})\left\{\frac{\partial}{\partial r}-\frac{1}{r}(\beta K-1)\right\}+\beta m+V(r)\right] \psi(\boldsymbol{r})=E \psi(\boldsymbol{r}) \tag{3.55}
\end{equation*}
$$

Since $\psi(\boldsymbol{r})$ is also an eigenfunction of the operator $K$ as shown in Eq. (3.43), we can substitute its eigenvalue $\kappa$ in Eq. (3.55).

$$
\begin{equation*}
\left[-i(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})\left\{\frac{\partial}{\partial r}-\frac{1}{r}(\beta \kappa-1)\right\}+\beta m+V(r)\right] \psi(\boldsymbol{r})=E \psi(\boldsymbol{r}) . \tag{3.56}
\end{equation*}
$$

[^18]Equation (3.56) is a matrix equation that can be written explicitly as

$$
\left[\begin{array}{cc}
m+V(\boldsymbol{r}) & X_{1}  \tag{3.57}\\
X_{2} & -m+V(\boldsymbol{r})
\end{array}\right]\left[\begin{array}{c}
\Phi_{a}(\boldsymbol{r}) \\
\Phi_{b}(\boldsymbol{r})
\end{array}\right]=E\left[\begin{array}{c}
\Phi_{a}(\boldsymbol{r}) \\
\Phi_{b}(\boldsymbol{r})
\end{array}\right] .
$$

with

$$
X_{1}=-i(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})\left\{\frac{\partial}{\partial r}+\frac{1}{r}(\kappa+1)\right\} ; \quad X_{2}=-i(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})\left\{\frac{\partial}{\partial r}-\frac{1}{r}(\kappa-1)\right\} .
$$

Equation (3.57) is obtained using the two-component form (3.47) for $\psi(\boldsymbol{r})$ and it leads to two coupled equations in $\Phi_{a}(\boldsymbol{r})$ and $\Phi_{b}(\boldsymbol{r})$. This coupling is essentially due to the $\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}}$ term in Eq. (3.55).

$$
\begin{align*}
& (E-m-V(r)) \Phi_{a}(\boldsymbol{r})=-i(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})\left\{\frac{\partial}{\partial r}+\frac{1}{r}(\kappa+1)\right\} \Phi_{b}(\boldsymbol{r}) .  \tag{3.58}\\
& (E+m-V(r)) \Phi_{b}(\boldsymbol{r})=-i(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})\left\{\frac{\partial}{\partial r}-\frac{1}{r}(\kappa-1)\right\} \Phi_{a}(\boldsymbol{r}) . \tag{3.59}
\end{align*}
$$

### 3.3.2 Radial and Angular-spin functions

You may recall that in solving the Schrödinger equation for a spherically symmetric potential, we were able to write down the wave function as a product of radial function and spherical harmonics. That was possible since the orbital angular momentum is a constant of motion in the nonrelativistic case. In the present problem, the orbital angular momentum $\boldsymbol{L}$ is not conserved but only the total angular momentum $\boldsymbol{J}=\boldsymbol{L}+\boldsymbol{S}$ is conserved. So, let us write down the state functions $\Phi_{a}(\boldsymbol{r})$ and $\Phi_{b}(\boldsymbol{r})$ as a product of radial and angular-spin functions.

$$
\begin{align*}
\Phi_{a}(\boldsymbol{r}) & =R_{a}(r) \Omega_{\kappa_{a}, j, m}^{a}(\theta, \phi) ;  \tag{3.60}\\
\Phi_{b}(\boldsymbol{r}) & =R_{b}(r) \Omega_{\kappa_{b}, j, m}^{b}(\theta, \phi) . \tag{3.61}
\end{align*}
$$

For a given $j$, two orbital states $l_{a}=j-\frac{1}{2}$ and $l_{b}=j+\frac{1}{2}$ are possible and one can construct two angular-spin functions $\Omega_{\kappa_{a}, j, m}^{a}=\left|l_{a}, \frac{1}{2}, j, m\right\rangle$ and $\Omega_{\kappa_{b}, j, m}^{b}=\left|l_{b}, \frac{1}{2}, j, m\right\rangle$.

$$
\begin{align*}
\Omega_{\kappa_{a}, j, m}^{a}(\theta, \phi)= & {\left[\begin{array}{ccc}
l_{a} & \frac{1}{2} & j \\
m-\frac{1}{2} & \frac{1}{2} & m
\end{array}\right] Y_{l_{a}}^{m-\frac{1}{2}}(\theta, \phi) \chi_{\frac{1}{2}}^{\frac{1}{2}} } \\
& +\left[\begin{array}{ccc}
l_{a} & \frac{1}{2} & j \\
m+\frac{1}{2} & -\frac{1}{2} & m
\end{array}\right] Y_{l_{a}}^{m+\frac{1}{2}}(\theta, \phi) \chi_{\frac{1}{2}}^{-\frac{1}{2}} \tag{3.62}
\end{align*}
$$

Table 3.1: Clebsch-Gordan (C.G.) coefficients

$$
\left[\begin{array}{ccc}
l & \frac{1}{2} & j \\
m_{1} & m_{2} & m
\end{array}\right]
$$

| $j$ | $m_{2}=\frac{1}{2}$ | $m_{2}=-\frac{1}{2}$ |
| :---: | :---: | :---: |
| $l+\frac{1}{2}$ | $\left(\frac{l+m+\frac{1}{2}}{2 l+1}\right)^{\frac{1}{2}}$ | $\left(\frac{l-m+\frac{1}{2}}{2 l+1}\right)^{\frac{1}{2}}$ |
| $l-\frac{1}{2}$ | $-\left(\frac{l-m+\frac{1}{2}}{2 l+1}\right)^{\frac{1}{2}}$ | $\left(\frac{l+m+\frac{1}{2}}{2 l+1}\right)^{\frac{1}{2}}$ |

where the square brackets denote the C.G. coefficients ${ }^{6}, Y_{l}^{m}(\theta, \phi)$ denote the spherical harmonics and $\chi$ the Pauli spin vectors.

$$
\chi_{\frac{1}{2}}^{\frac{1}{2}}=\left[\begin{array}{l}
1 \\
0
\end{array}\right] ; \quad \chi_{\frac{1}{2}}^{-\frac{1}{2}}=\left[\begin{array}{l}
0 \\
1
\end{array}\right]
$$

Using the algebraic expressions for the C.G. coefficients (vide Table 3.1), we get

$$
\begin{align*}
\Omega_{\kappa_{a}, j, m}^{a}= & \left(\frac{l_{a}+m+\frac{1}{2}}{2 l_{a}+1}\right)^{1 / 2} Y_{l_{a}}^{m-\frac{1}{2}}\left[\begin{array}{l}
1 \\
0
\end{array}\right] \\
& +\left(\frac{l_{a}-m+\frac{1}{2}}{2 l_{a}+1}\right)^{1 / 2} Y_{l_{a}}^{m+\frac{1}{2}}\left[\begin{array}{l}
0 \\
1
\end{array}\right] \\
= & {\left[\begin{array}{l}
\left(\frac{l_{a}+m+\frac{1}{2}}{2 l_{a}+1}\right)^{1 / 2} Y_{l_{a}}^{m-\frac{1}{2}} \\
\left(\frac{l_{a}-m+\frac{1}{2}}{2 l_{a}+1}\right)^{1 / 2} Y_{l_{a}}^{m+\frac{1}{2}}
\end{array}\right] . } \tag{3.63}
\end{align*}
$$

In a similar way, the other angular spin-function can be written as

$$
\Omega_{\kappa_{b}, j, m}^{b}=\left[\begin{array}{c}
-\left(\frac{l_{b}-m+\frac{1}{2}}{2 l_{b}+1}\right)^{1 / 2} Y_{l_{b}}^{m-\frac{1}{2}}  \tag{3.64}\\
\left(\frac{l_{b}+m+\frac{1}{2}}{2 l_{b}+1}\right)^{1 / 2} Y_{l_{b}}^{m+\frac{1}{2}}
\end{array}\right]
$$

[^19]The parity of the two angular-spin functions are given by $(-1)^{l_{a}}$ and $(-1)^{l_{b}}$, and so they are of opposite parity since $\left|l_{a}-l_{b}\right|=1$. It can be easily verified (vide Solved Problem 2.5) that

$$
\begin{align*}
\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} \Omega_{\kappa_{a}, j, m}^{a} & =-\Omega_{\kappa_{b}, j, m}^{b}  \tag{3.65}\\
\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} \Omega_{\kappa_{b}, j, m}^{b} & =-\Omega_{\kappa_{a}, j, m}^{a} \tag{3.66}
\end{align*}
$$

The operator $K$ has two eigenvalues $\kappa_{a}$ and $\kappa_{b}$ as given by Eq. (3.46).

$$
\kappa_{a, b}=\mp\left(j+\frac{1}{2}\right) .
$$

It is customary to take $\kappa_{a}=-\left(j+\frac{1}{2}\right)$ and $\kappa_{b}=+\left(j+\frac{1}{2}\right)$. It can be easily verified that

$$
\begin{array}{lll}
\kappa_{a}=-\left(j+\frac{1}{2}\right) ; & j=\ell_{a}+\frac{1}{2} ; & \kappa_{a}=-\left(\ell_{a}+1\right) . \\
\kappa_{b}=+\left(j+\frac{1}{2}\right) ; & j=\ell_{b}-\frac{1}{2} ; & \kappa_{b}=\ell_{b} . \tag{3.68}
\end{array}
$$

Note that $\left|\kappa_{a}\right|=\left|\kappa_{b}\right|=|\kappa|$.

### 3.3.3 Coupled Equations for Radial functions

It is remarkable that the pseudo-scalar operator $\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}}$ switches the one angular-spin function to the other and thereby enables you to separate the radial equations from the angular dependent quantities. In the coupled equations (3.58) and (3.59), the pseudo-scalar operator $\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}}$ does precisely switch the angular-spin functions. Substituting (3.60) and (3.61) for $\Phi_{a}$ and $\Phi_{b}$ and using Eqs. (3.65) and (3.66) in Eqs. (3.58) and (3.59) and eliminating the same angular-spin functions that occur on both sides, we obtain the following coupled equations for the radial functions $R_{a}(r)$ and $R_{b}(r)$.

$$
\begin{align*}
& \{E-m-V(r)\} R_{a}(r)=i\left\{\frac{\partial}{\partial r}+\frac{1}{r}(\kappa+1)\right\} R_{b}(r) .  \tag{3.69}\\
& \{E+m-V(r)\} R_{b}(r)=i\left\{\frac{\partial}{\partial r}-\frac{1}{r}(\kappa-1)\right\} R_{a}(r) . \tag{3.70}
\end{align*}
$$

Let us write the radial functions as

$$
\left[\begin{array}{c}
R_{a}  \tag{3.71}\\
R_{b}
\end{array}\right]=\left[\begin{array}{c}
(1 / r) F(r) \\
(i / r) G(r)
\end{array}\right],
$$

where the phase factor $i$ is introduced for the function $G(r)$ to make the radial equations for $F(r)$ and $G(r)$ explicitly real. Since

$$
\begin{align*}
\frac{\partial}{\partial r}\left(\frac{F(r)}{r}\right) & =\frac{1}{r} \frac{\partial F(r)}{\partial r}-\frac{1}{r^{2}} F(r)  \tag{3.72}\\
\frac{\partial}{\partial r}\left(\frac{G(r)}{r}\right) & =\frac{1}{r} \frac{\partial G(r)}{\partial r}-\frac{1}{r^{2}} G(r) \tag{3.73}
\end{align*}
$$

the coupled equations for radial functions $F(r)$ and $G(r)$ become

$$
\begin{align*}
& \{E-m-V(r)\} F(r)+\frac{\partial G}{\partial r}+\frac{\kappa}{r} G(r)=0 ;  \tag{3.74}\\
& \{E+m-V(r)\} G(r)-\frac{\partial F}{\partial r}+\frac{\kappa}{r} F(r)=0 \tag{3.75}
\end{align*}
$$

Let

$$
\begin{equation*}
a_{1}=m+E, \quad a_{2}=m-E, \quad a=\sqrt{a_{1} a_{2}}=\sqrt{m^{2}-E^{2}} . \tag{3.76}
\end{equation*}
$$

and introduce a new variable

$$
\begin{equation*}
\rho=a r, \quad \text { such that } \quad \frac{\partial}{\partial \rho} \frac{\partial \rho}{\partial r}=a \frac{\partial}{\partial \rho} . \tag{3.77}
\end{equation*}
$$

Let us rewrite the radial equations (3.74) and (3.75) in terms of the new variable $\rho$.

$$
\begin{align*}
& \left(\frac{\partial}{\partial \rho}+\frac{\kappa}{\rho}\right) G(\rho)-\frac{1}{a}\left\{a_{2}+V(r)\right\} F(\rho)=0  \tag{3.78}\\
& \left(\frac{\partial}{\partial \rho}-\frac{\kappa}{\rho}\right) F(\rho)-\frac{1}{a}\left\{a_{2}-V(r)\right\} G(\rho)=0 \tag{3.79}
\end{align*}
$$

### 3.4 Hydrogen-like atom

For hydrogen-like atom, the electron experiences the Coulomb potential due to the nucleus of charge $Z e$.

$$
\begin{equation*}
V(r)=-\frac{Z e^{2}}{r}=-a \frac{Z e^{2}}{\rho}=-\frac{a \Gamma}{\rho}, \tag{3.80}
\end{equation*}
$$

where $\Gamma=Z e^{2}=Z \alpha$. The symbol $\alpha$ denotes the fine structure constant ( $\alpha=1 / 137$ ). For finding the energy levels of the hydrogen-like atom, we
need to solve the radial equations with the above Coulomb potential.

$$
\begin{align*}
& \left(\frac{\partial}{\partial \rho}+\frac{\kappa}{\rho}\right) G(\rho)-\frac{1}{a}\left(a_{2}-\frac{a \Gamma}{\rho}\right) F(\rho)=0  \tag{3.81}\\
& \left(\frac{\partial}{\partial \rho}-\frac{\kappa}{\rho}\right) F(\rho)-\frac{1}{a}\left(a_{1}+\frac{a \Gamma}{\rho}\right) G(\rho)=0 . \tag{3.82}
\end{align*}
$$

Since we are interested in the bound state solutions, the radial functions should vanish at infinity. So, we shall choose the radial functions to be of the form

$$
\begin{equation*}
G(\rho)=g(\rho) e^{-\rho} ; \quad F(\rho)=f(\rho) e^{-\rho} . \tag{3.83}
\end{equation*}
$$

Then

$$
\begin{equation*}
\frac{\partial}{\partial \rho} G(\rho)=\frac{\partial}{\partial \rho}\left(g(\rho) e^{-\rho}\right)=\left(\frac{\partial g}{\partial \rho}-g\right) e^{-\rho}=\left(g^{\prime}-g\right) e^{-\rho} . \tag{3.84}
\end{equation*}
$$

Similarly,

$$
\begin{equation*}
\frac{\partial}{\partial \rho} F(\rho)=\frac{\partial}{\partial \rho}\left(f(\rho) e^{-\rho}\right)=\left(\frac{\partial f}{\partial \rho}-f\right) e^{-\rho}=\left(f^{\prime}-f\right) e^{-\rho} \tag{3.85}
\end{equation*}
$$

The radial equations (3.81) and (3.82) in terms of $f(\rho)$ and $g(\rho)$ become

$$
\begin{align*}
& g^{\prime}-g+\frac{\kappa}{\rho} g-\left(\frac{a_{2}}{a}-\frac{\Gamma}{\rho}\right) f=0  \tag{3.86}\\
& f^{\prime}-f-\frac{\kappa}{\rho} f-\left(\frac{a_{1}}{a}+\frac{\Gamma}{\rho}\right) g=0 . \tag{3.87}
\end{align*}
$$

### 3.4.1 Power series expansion for radial functions

We shall adopt the usual technique of power series expansion for $f(\rho)$ and $g(\rho)$ and examine their behaviour for large $\rho$. If they diverge, then the series should be terminated with a finite number of terms to satisfy the asymptotic boundary condition of radial functions $F(\rho)$ and $G(\rho)$, vanishing at large $\rho$.

Let us make the power series expansion for $f(\rho)$ and $g(\rho)$.

$$
\begin{align*}
& f(\rho)=\rho^{s}\left(b_{0}+b_{1} \rho+b_{2} \rho^{2}+\cdots\right)=\rho^{s} \sum_{\nu=0}^{\infty} b_{\nu} \rho^{\nu}, \quad b_{0} \neq 0 .  \tag{3.88}\\
& g(\rho)=\rho^{s}\left(c_{0}+c_{1} \rho+c_{2} \rho^{2}+\cdots\right)=\rho^{s} \sum_{\nu=0}^{\infty} c_{\nu} \rho^{\nu}, \quad c_{0} \neq 0 . \tag{3.89}
\end{align*}
$$

Then

$$
\begin{align*}
& f^{\prime}(\rho)=\sum_{\nu=0}^{\infty}(s+\nu) b_{\nu} \rho^{s+\nu-1}  \tag{3.90}\\
& g^{\prime}(\rho)=\sum_{\nu=0}^{\infty}(s+\nu) c_{\nu} \rho^{s+\nu-1} \tag{3.91}
\end{align*}
$$

Substituting the power series expansions (3.88) - (3.91) into Eqs. (3.86) and (3.87), let us find the coefficient of the term $\rho^{s+\nu-1}$.

$$
\begin{align*}
& (s+\nu+\kappa) c_{\nu}-c_{\nu-1}-\frac{a_{2}}{a} b_{\nu-1}+\Gamma b_{\nu}=0  \tag{3.92}\\
& (s+\nu-\kappa) b_{\nu}-b_{\nu-1}-\frac{a_{1}}{a} c_{\nu-1}-\Gamma c_{\nu}=0 . \tag{3.93}
\end{align*}
$$

Since the negative subscripts are not allowed for $b$ and $c$, the substitution $\nu=0$ in Eqs. (3.92) and (3.93) yield

$$
\begin{align*}
& (s+\kappa) c_{0}+\Gamma b_{0}=0  \tag{3.94}\\
& (s-\kappa) b_{0}-\Gamma c_{0}=0 \tag{3.95}
\end{align*}
$$

which can be written in the form of a matrix equation

$$
\left[\begin{array}{cc}
s-\kappa & -\Gamma  \tag{3.96}\\
\Gamma & s+\kappa
\end{array}\right]\left[\begin{array}{l}
b_{0} \\
c_{0}
\end{array}\right]=0 .
$$

Since $b_{0}$ and $c_{0}$ are not zero, to obtain a non-trivial solution, the condition is that the determinant of the secular matrix should be zero.

$$
\begin{equation*}
s^{2}-\kappa^{2}+\Gamma^{2}=0 \quad \text { or } \quad s=+\sqrt{\kappa^{2}-\Gamma^{2}}, \tag{3.97}
\end{equation*}
$$

since $s$ is a positive quantity. This implies that $|\kappa|<\Gamma=Z \alpha$ or $Z<\frac{|\kappa|}{\alpha}$. This means that there is no solution for one-electron (hydrogen-like) atom if $Z>137$, since the minimum value of $|\kappa|=1$.

In order to eliminate the terms $b_{\nu-1}$ and $c_{\nu-1}$ from Eqs. (3.92) and (3.93), multiply Eq. (3.92) by $a_{1}$ and Eq. (3.93) by $a$ and subtract.

$$
\begin{align*}
a_{1}(s+\nu+\kappa) c_{\nu}-a_{1} c_{\nu-1}-\frac{a_{1} a_{2}}{a} b_{\nu-1}+a_{1} \Gamma b_{\nu} & =0 ;  \tag{3.98}\\
a(s+\nu-\kappa) b_{\nu}-a b_{\nu-1}-a_{1} c_{\nu-1}-a \Gamma c_{\nu} & =0 . \tag{3.99}
\end{align*}
$$

Subtracting (3.99) from (3.98) and recalling that $a_{1} a_{2}=a^{2}$, we get

$$
\begin{equation*}
c_{\nu}\left\{a_{1}(s+\nu+\kappa)+a \Gamma\right\}=b_{\nu}\left\{a(s+\nu-\kappa)-a_{1} \Gamma\right\} . \tag{3.100}
\end{equation*}
$$

For bound states, the radial wave functions must be finite near the origin but should vanish for large $r$. If $s \geq 1$, the wave function will be finite near the origin. For larger $r$, the wave function will vanish only if $F(r)$ and $G(r)$ are finite at large $r$. The behaviour at large $r$ is determined essentially by large $\nu$. Omitting the other terms and retaining only the $\nu$-dependent terms, Eqs. (3.92) and (3.93) can be written as

$$
\begin{align*}
& \nu c_{\nu}-c_{\nu-1}-\frac{a_{2}}{a} b_{\nu-1}+\Gamma b_{\nu}=0 ;  \tag{3.101}\\
& \nu b_{\nu}-b_{\nu-1}-\frac{a_{1}}{a} c_{\nu-1}-\Gamma c_{\nu}=0 . \tag{3.102}
\end{align*}
$$

For large $\nu$, Eq. (3.100) yields a relation between the coefficients $b_{\nu}$ and $c_{\nu}$.

$$
\begin{equation*}
a_{1} c_{\nu}=a b_{\nu} \tag{3.103}
\end{equation*}
$$

Using the relation (3.103) in Eqs. (3.101) and (3.102), we get

$$
\begin{equation*}
\Gamma=\frac{\nu b_{\nu}-2 b_{\nu-1}}{c_{\nu}}=\frac{2 c_{\nu-1}-\nu c_{\nu}}{b_{\nu}} . \tag{3.104}
\end{equation*}
$$

Equivalently

$$
\begin{equation*}
b_{\nu}\left(\nu b_{\nu}-2 b_{\nu-1}\right)+c_{\nu}\left(\nu c_{\nu}-2 c_{\nu-1}\right)=0 . \tag{3.105}
\end{equation*}
$$

Since the coefficients $b$ and $c$ are independent, the above relation is true only if each term vanishes separately for large $r$.

$$
\begin{equation*}
b_{\nu} \rightarrow \frac{2 b_{\nu-1}}{\nu} \quad \text { and } \quad c_{\nu} \rightarrow \frac{2 c_{\nu-1}}{\nu} . \tag{3.106}
\end{equation*}
$$

Thus for large $r$,

$$
\begin{array}{ll}
b_{1}=2 b_{0} & c_{1}=2 c_{0} \\
b_{2}=\frac{2 b_{1}}{2}=2 b_{0} & c_{2}=c_{1}=2 c_{0} \\
b_{3}=\frac{2 b_{2}}{3}=\frac{4}{3} b_{0} & c_{3}=\frac{2 c_{2}}{3}=\frac{4}{3} c_{0}  \tag{3.107}\\
b_{4}=\frac{2 b_{3}}{4}=\frac{2}{3} b_{0} & c_{4}=\frac{2 c_{3}}{4}=\frac{2}{3} c_{0}
\end{array}
$$

Substituting these coefficients $b_{\nu}$ and $c_{\nu}$ in Eqs. (3.88) and (3.89), we
obtain

$$
\begin{align*}
f(\rho) & =\rho^{s} \sum_{\nu=0}^{\infty} b_{\nu} \rho^{\nu} \\
& =\rho^{s} b_{0}\left\{1+2 \rho+2 \rho^{2}+\frac{4}{3} \rho^{3}+\frac{2}{3} \rho^{4}+\cdots\right\} \\
& =\rho^{s} b_{0}\left\{1+2 \rho+\frac{1}{2!}(2 \rho)^{2}+\frac{1}{3!}(2 \rho)^{3}+\frac{1}{4!}(2 \rho)^{4}+\cdots\right\} \\
& =b_{0} \rho^{s} e^{2 \rho} . \tag{3.108}
\end{align*}
$$

Similarly,

$$
\begin{equation*}
g(\rho)=c_{0} \rho^{s} e^{2 \rho} \tag{3.109}
\end{equation*}
$$

Since the series diverges as $e^{2 \rho}$, in order to obtain a well-behaved function, the series has to be terminated or truncated at some value of $\nu$, say $\nu=n_{r}$, so that $b_{n_{r}+1}=c_{n_{r}+1}=0$. From Eqs. (3.92) and (3.93), we obtain a relationship between $b_{n_{r}}$ and $c_{n_{r}}$ by substituting $\nu-1=n_{r}$.

$$
\begin{equation*}
a c_{n_{r}}=-a_{2} b_{n_{r}} ; \quad a b_{n_{r}}=-a_{1} c_{n_{r}} . \tag{3.110}
\end{equation*}
$$

Substituting $\nu=n_{r}$ in Eq. (3.100), we get

$$
\begin{equation*}
c_{n_{r}}\left\{a_{1}\left(s+n_{r}+\kappa\right)+a \Gamma\right\}=b_{n_{r}}\left\{a\left(s+n_{r}-\kappa\right)-a_{1} \Gamma\right\} . \tag{3.111}
\end{equation*}
$$

By using the relation (3.110) in Eq. (3.111), the coefficients $b$ and $c$ are eliminated and an expression for the energy levels of the hydrogen-like atom is obtained.

$$
\begin{equation*}
\left(a_{1}-a_{2}\right) \Gamma=2 a\left(s+n_{r}\right) \quad \text { or } \quad 2 E \Gamma \quad=2 a\left(s+n_{r}\right), \tag{3.112}
\end{equation*}
$$

since $a_{1}=m+E$ and $a_{2}=m-E$. Also $a^{2}=a_{1} a_{2}=m^{2}-E^{2}$. Squaring (3.112), we obtain

$$
\begin{align*}
4 E^{2} \Gamma^{2} & =4 a^{2}\left(s+n_{r}\right)^{2}=4\left(m^{2}-E^{2}\right)\left(s+n_{r}\right)^{2} ; \\
E^{2} & =\frac{m^{2}\left(s+n_{r}\right)^{2}}{\left(s+n_{r}\right)^{2}+\Gamma^{2}}=\frac{m^{2}}{1+\frac{\Gamma^{2}}{\left(s+n_{r}\right)^{2}}} . \tag{3.113}
\end{align*}
$$

### 3.4.2 Energy levels

Taking the square root of (3.113), we obtain an expression for the energy levels of hydrogen-like atom.

$$
\begin{equation*}
E=m\left\{1+\frac{\Gamma^{2}}{\left(s+n_{r}\right)^{2}}\right\}^{-\frac{1}{2}} \tag{3.114}
\end{equation*}
$$

This is an exact relativistic expression for energy but an approximate expression in powers of $\alpha^{2}$ or $\Gamma^{2}$ can be obtained using binomial expansion.

The symbol $s$ is defined by Eq. (3.97).

$$
\begin{align*}
s & =\sqrt{\kappa^{2}-\Gamma^{2}}=|\kappa|\left(1-\frac{\Gamma^{2}}{\kappa^{2}}\right)^{\frac{1}{2}} \approx|\kappa|\left(1-\frac{\Gamma^{2}}{2 \kappa^{2}}\right) . \\
s+n_{r} & =|\kappa|+n_{r}-\frac{\Gamma^{2}}{2|\kappa|}=n-\frac{\Gamma^{2}}{2|\kappa|}, \tag{3.115}
\end{align*}
$$

where $n=|\kappa|+n_{r}$ is the principal quantum number and $n_{r}$ is known as the radial quantum number. The energy levels of the hydrogen-like atom is given by

$$
\begin{align*}
E & =m\left\{1+\frac{\Gamma^{2}}{\left(n-\frac{\Gamma^{2}}{2|\kappa|}\right)^{2}}\right\}^{-\frac{1}{2}} \\
& =m\left\{1-\frac{1}{2} \frac{\Gamma^{2}}{\left(n-\frac{\Gamma^{2}}{2|\kappa|}\right)^{2}}+\frac{\left(-\frac{1}{2}\right)\left(-\frac{3}{2}\right)}{2!} \frac{\Gamma^{4}}{\left(n-\frac{\Gamma^{2}}{2|\kappa|}\right)^{4}}+\cdots\right\} \\
& =m\left\{1-\frac{\Gamma^{2}}{2 n^{2}}\left(1-\frac{\Gamma^{2}}{2 n|\kappa|}\right)^{-2}+\frac{3 \Gamma^{4}}{8 n^{4}}\left(1-\frac{\Gamma^{2}}{2 n|\kappa|}\right)^{-4}+\cdots\right\} \\
& =m\left\{1-\frac{\Gamma^{2}}{2 n^{2}}\left(1+\frac{\Gamma^{2}}{n|\kappa|}\right)+\frac{3 \Gamma^{4}}{8 n^{4}}+\cdots\right\} \\
& =m\left\{1-\frac{\Gamma^{2}}{2 n^{2}}-\frac{\Gamma^{4}}{2 n^{4}}\left(\frac{n}{|\kappa|}-\frac{3}{4}\right)+\cdots\right\} . \tag{3.116}
\end{align*}
$$

Let us examine the expression (3.116) obtained for the energy level in natural units $\hbar=c=1$. The first term $m c^{2}$ (when restored to C.G.S. or M.K.S. units) denotes the rest energy of the electron, the second term the non-relativistic Bohr energy and the third term corresponds to the fine
structure effect depending on the quantum number $\kappa$. The energy depends only on $|\kappa|$ and not on its sign. Since $\kappa$ can assume both positive and negative values $\pm 1, \pm 2, \pm 3, \cdots$, there is a degeneracy in the energy levels. The following pairs of energy levels with the same $|\kappa|$ are degenerate.

$$
\left(2 S_{1 / 2}, 2 P_{1 / 2}\right),\left(3 S_{1 / 2}, 3 P_{1 / 2}\right),\left(3 P_{3 / 2}, 3 D_{3 / 2}\right), \cdots
$$

The degenerate pairs of energy levels correspond to two eigenstates of opposite parity corresponding to the same quantum numbers $n$ and $j$.

The energy levels are given in Table 3.2 and depicted in Fig. 3.1. Besides using the relativistic quantum numbers, $n, \kappa, j, m$, the energy levels are also designated, by convention, by spectroscopic notation $n\{\ell\}_{j}$, using the non-relativistic concept of spin-orbit interaction.

Table 3.2: Low-lying energy levels of hydrogen atom, characterized by relativistic quantum numbers $n, \kappa, j$. The conventional spectroscopic notation $n\{\ell\}_{j}$ is also included.

| $n$ | $n_{r}$ | $\kappa$ | $\ell$ | $j$ | Spectroscopic <br> notation $n\{\ell\}_{j}$ | Energy level <br> $E_{n \kappa j}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 0 | -1 | 0 | $\frac{1}{2}$ | $1 \mathrm{~S}_{\frac{1}{2}}$ | $m\left(1-\frac{\alpha^{2}}{2}-\frac{\alpha^{4}}{8}\right)$ |
| 2 | 1 | -1 | 0 | $\frac{1}{2}$ | $2 \mathrm{~S}_{\frac{1}{2}}$ | $m\left(1-\frac{\alpha^{2}}{8}-\frac{5 \alpha^{4}}{128}\right)$ |
|  | 1 | +1 | 1 | $\frac{1}{2}$ | $2 \mathrm{P}_{\frac{1}{2}}^{2}$ | $m\left(1-\frac{\alpha^{2}}{8}-\frac{\alpha^{4}}{128}\right)$ |
|  | 0 | -2 | 1 | $\frac{3}{2}$ | $2 \mathrm{P}_{\frac{3}{2}}$ | $m\left(1-\frac{\alpha^{2}}{18}-\frac{9 \alpha^{4}}{648}\right)$ |
| 3 | 2 | -1 | 0 | $\frac{1}{2}$ | $3 \mathrm{~S}_{\frac{1}{2}}$ | $m m\left(1-\frac{\alpha^{2}}{18}-\frac{3 \alpha^{4}}{648}\right)$ |
|  | 2 | +1 | 1 | $\frac{1}{2}$ | $3 \mathrm{P}_{\frac{1}{2}}$ |  |
|  | 1 | -2 | 1 | $\frac{3}{2}$ | $3 \mathrm{P}_{\frac{3}{2}}$ | $m m\left(1-\frac{\alpha^{2}}{18}-\frac{\alpha^{4}}{648}\right)$ |

From the Table 3.2, it can be observed that for the same $n$, there is a splitting of energy levels due to different $|\kappa|$ values and this is known as

(a)

$$
\frac{2 S_{1 / 2}}{2 P_{1 / 2}}{\frac{2 S_{1 / 2}}{2 P_{1 / 2}}}_{\overbrace{(1057 \mathrm{MHz})}}^{\substack{E=4.372 \times 10^{-6} \mathrm{eV} \\(10 y)}}
$$

Lamb Shift

(b)

Figure 3.1: (a) Low-lying energy levels of hydrogen atom as per Dirac's theory. The $S, P, D$ states corresponding to $\ell=0,1,2$ are shown laterally displaced for the sake of clarity. One can observe the pairs of degenerate states with the same $j$ but different $\ell$ values. (b) The hyperfine splitting for $1 S_{1 / 2}$ level and the Lamb shift for $2 S_{1 / 2}-2 P_{1 / 2}$ level are shown. The diagrams are not drawn to scale.
the fine structure splitting due to spin-orbit interaction.

$$
\begin{align*}
& E\left(2 P_{3 / 2}\right)-E\left(2 P_{1 / 2}\right)=\frac{m \alpha^{4}}{32} \approx 4.5 \times 10^{-5} \mathrm{eV} .(\nu=10.9 \mathrm{GHz} .)  \tag{3.117}\\
& E\left(3 P_{3 / 2}\right)-E\left(3 P_{1 / 2}\right)=\frac{m \alpha^{4}}{108} \approx 1.33 \times 10^{-5} \mathrm{eV} .(\nu=3.22 \mathrm{GHz} .)  \tag{3.118}\\
& E\left(3 D_{5 / 2}\right)-E\left(3 D_{3 / 2}\right)=\frac{m \alpha^{4}}{324} \approx 4.44 \times 10^{-6} \mathrm{eV} .(\nu=1.075 \mathrm{GHz} .) \tag{3.119}
\end{align*}
$$

The experimental observations completely confirm the predictions of Dirac's theory and the fine structure of hydrogen spectral lines arising from spin-orbit interaction. However, it is necessary to include the inter-
action of electron and nuclear (proton) spins. This gives rise to hyperfine structure of spectral lines, each energy level being split into two levels corresponding to the two possible total angular momentum states $F$ that can be obtained by coupling of $j$ of electron with the proton spin in hydrogen atom. The hyperfine splitting of $1 S_{1 / 2}$ level is found to be $\Delta E=5.9 \times 10^{-6}$ eV that corresponds to a frequency of 1420 MHz . The hyperfine splitting of energy level $1 S_{1 / 2}$ into two is alone shown in Fig. 3.1(b) but such a splitting occurs for each level. So, each level becomes a doublet due to hyperfine interaction.

Further, in 1947, Lamb-Retherford ${ }^{7}$ found that the degeneracy of energy levels with the same $n$ and $j$ but differing $\ell$ is removed due to the interaction of electrons with the fluctuations of the quantized radiation field as envisaged in Quantum Field Theory. It is found that the $2 S_{1 / 2}$ energy level is slightly higher than the $2 P_{1 / 2}$ and this difference in energy is known as the Lamb shift.

$$
E\left(2 S_{1 / 2}\right)-E\left(2 P_{1 / 2}\right)=4.372 \times 10^{-6} \mathrm{eV}=1057 \mathrm{MHz}
$$

Both the hyperfine splitting and the Lamb shift have been measured and are in perfect agreement with the theoretical calculations. This is considered as one of the greatest triumphs of Quantum Electrodynamics.

## Review Questions

3.1 Write down the Dirac equation for a charged particle in electromagnetic field and show that its non-relativistic reduction automatically yields terms which represent magnetic moment interaction with the correct gyromagnetic ratio for the spin, spin-orbital interaction, Darwin term and Thomas term, which have been studied earlier with non-relativistic Schrödinger equation with the inclusion of Pauli spin term in an ad hoc manner.
3.2 Write down the Dirac equation for a central potential and show how you can separate the radial and angular spin functions in the eigenstates to obtain a set of two coupled radial wave equations.
3.3 Explain how you can solve the coupled radial wave equations obtained for the Dirac equation with Coulomb potential with the bound state boundary conditions and obtain the energy levels of hydrogen-like atom. Show that the Dirac theory yields correctly the fine structure of spectral lines arising from spin-orbit interaction.
3.4 It is found experimentally that the hydrogen spectrum shows not only the fine structure predicted by Dirac's theory but also a hyperfine structure.

[^20]Besides it is observed that the degeneracy of energy levels with the same $j$ but different $\ell$ as predicted by Dirac's theory is not true and a small splitting is observed. How will you account for these discrepancies?

## Problems

3.1 Write down explicitly the operator $K=\beta(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1)$ in the matrix form and show that it is Hermitian.
3.2 Show that the Dirac Hamiltonian $H$ in a central potential commutes with the following operators

$$
K=\beta(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1), \quad \text { and } \quad \boldsymbol{J}=\boldsymbol{L}+\boldsymbol{S}
$$

and that the operators $K$ and $\boldsymbol{J}$ commute among themselves.
3.3 Deduce the relation (3.52)

$$
\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}} \times \boldsymbol{L}=-i(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \boldsymbol{L})
$$

3.4 Prove that

$$
(\alpha \cdot \hat{\boldsymbol{r}})^{2}=1
$$

Using this relation, derive Eq. (3.54)

$$
\boldsymbol{\alpha} \cdot \boldsymbol{p}=-i(\alpha \cdot \hat{\boldsymbol{r}})\left\{\frac{\partial}{\partial r}-\frac{1}{r}(\beta K-1)\right\}
$$

3.5 Prove that the matrix element

$$
Q=\left\langle l_{f} \frac{1}{2} j_{f} m_{f}\right| \boldsymbol{\sigma} \cdot \boldsymbol{r}\left|l_{i} \frac{1}{2} j_{i} m_{i}\right\rangle=-\delta_{j_{i} j_{f}} \delta_{m_{i} m_{f}}
$$

and hence show that

$$
\begin{aligned}
\boldsymbol{\sigma} \cdot \boldsymbol{r}\left|l_{a} \frac{1}{2} j m\right\rangle & =-\left|l_{b} \frac{1}{2} j m\right\rangle \\
\boldsymbol{\sigma} \cdot \boldsymbol{r}\left|l_{b} \frac{1}{2} j m\right\rangle & =-\left|l_{a} \frac{1}{2} j m\right\rangle
\end{aligned}
$$

if $l_{a}=j-\frac{1}{2}$ and $l_{b}=j+\frac{1}{2}$.
3.6 (a) Obtain the following relations:
(i) $\boldsymbol{L}^{2}=\boldsymbol{r}^{2} \boldsymbol{p}^{2}-(\boldsymbol{r} \cdot \boldsymbol{p})^{2}+i \hbar(\boldsymbol{r} \cdot \boldsymbol{p})$.
(ii) $p_{r}=\frac{1}{r}(\boldsymbol{r} \cdot \boldsymbol{p}-i \hbar)$.

The symbols $\boldsymbol{L}, \boldsymbol{r}, \boldsymbol{p}, p_{r}$ denote respectively, the orbital angular momentum, position vector, momentum vector, radial component of momentum in central field.
(b) Show that $p_{r}$ satisfies the commutation relation

$$
\left[r, p_{r}\right]_{-}=i \hbar
$$

(c) Prove that $p_{r}$ is a Hermitian operator.
3.7 Given the Dirac Hamiltonian $H=\boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m+V(r)$ in a central potential, express it in the form

$$
H=\alpha_{r} p_{r}+\frac{i}{r} \alpha_{r} \beta K+\beta m+V(r)
$$

using the operators $\alpha_{r}, p_{r}, K$, defined by

$$
\alpha_{r}=\frac{1}{r}(\boldsymbol{\alpha} \cdot \boldsymbol{r}) ; \quad p_{r}=\frac{1}{r}(\boldsymbol{r} \cdot \boldsymbol{p}-i) ; \quad K=\beta(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1)
$$

## Solutions to Problems

3.1 Let us write down explicitly the operator $K$ in the matrix form.

$$
\begin{aligned}
K & =\beta(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1)=\left[\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right]\left[\begin{array}{cc}
\boldsymbol{\sigma} \cdot \boldsymbol{L}+1 & 0 \\
0 & \boldsymbol{\sigma} \cdot \boldsymbol{L}+1
\end{array}\right] \\
& =\left[\begin{array}{cc}
\boldsymbol{\sigma} \cdot \boldsymbol{L}+1 & 0 \\
0 & -(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1)
\end{array}\right]
\end{aligned}
$$

Let us now expand $\boldsymbol{\sigma} \cdot \boldsymbol{L}$ and write it in the matrix form

$$
\begin{aligned}
\boldsymbol{\sigma} \cdot \boldsymbol{L} & =\sigma_{x} L_{x}+\sigma_{y} L_{y}+\sigma_{z} L_{z} \\
& =\left[\begin{array}{cc}
0 & L_{x} \\
L_{x} & 0
\end{array}\right]+\left[\begin{array}{cc}
0 & -i L_{y} \\
i L_{y} & 0
\end{array}\right]+\left[\begin{array}{cc}
L_{z} & 0 \\
0 & -L_{z}
\end{array}\right] \\
& =\left[\begin{array}{cc}
L_{z} & L_{x}-i L_{y} \\
L_{x}+i L_{y} & -L z
\end{array}\right] .
\end{aligned}
$$

It follows that $K$ can be written as a $4 \times 4$ matrix.

$$
\begin{aligned}
K & =\left[\begin{array}{ccc}
\boldsymbol{\sigma} \cdot \boldsymbol{L}+1 & 0 \\
0 & -(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1)
\end{array}\right] \\
& =\left[\begin{array}{cccc}
L_{z}+1 & L_{x}-i L_{y} & 0 & 0 \\
L_{x}+i L_{y} & -L_{z}+1 & 0 & 0 \\
0 & 0 & -L_{z}-1 & -L_{z}+i L_{y} \\
0 & 0 & -L_{x}-i L_{y} & L_{z}-1
\end{array}\right]
\end{aligned}
$$

The above matrix is Hermitian and hence $K^{\dagger}=K$.
3.2 The Dirac Hamiltonian $H$ and the operator $K$ are given by

$$
H=\boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m+V(r) ; \quad K=\beta(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1)
$$

## To show that $[H, K]_{-}=0$

The potential $V(r)$ does not involve any operator and so it commutes with the operator $K$. The mass term also commutes with $K$ since $\beta^{2}=1$ and $[\beta, \boldsymbol{\sigma} \cdot \boldsymbol{L}]_{-}=0$.

$$
[\beta, \beta(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1)]_{-}=0
$$

Only we need to evaluate the commutator $[\boldsymbol{\alpha} \cdot \boldsymbol{p}, K]_{-}$.

$$
\begin{align*}
{[\boldsymbol{\alpha} \cdot \boldsymbol{p}, \beta \boldsymbol{\sigma} \cdot \boldsymbol{L}]_{-} } & =(\boldsymbol{\alpha} \cdot \boldsymbol{p}) \beta(\boldsymbol{\sigma} \cdot \boldsymbol{L})-\beta(\boldsymbol{\sigma} \cdot \boldsymbol{L})(\boldsymbol{\alpha} \cdot \boldsymbol{p}) \\
& =-\beta(\boldsymbol{\alpha} \cdot \boldsymbol{p})(\boldsymbol{\sigma} \cdot \boldsymbol{L})-\beta(\boldsymbol{\sigma} \cdot \boldsymbol{L})(\boldsymbol{\alpha} \cdot \boldsymbol{p}) \\
& =-\beta\{(\boldsymbol{\alpha} \cdot \boldsymbol{p})(\boldsymbol{\sigma} \cdot \boldsymbol{L})+(\boldsymbol{\sigma} \cdot \boldsymbol{L})(\boldsymbol{\alpha} \cdot \boldsymbol{p})\} \cdot  \tag{3.120}\\
(\boldsymbol{\alpha} \cdot \boldsymbol{p})(\boldsymbol{\sigma} \cdot \boldsymbol{L}) & =\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \boldsymbol{p} \\
\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0
\end{array}\right]\left[\begin{array}{cc}
\boldsymbol{\sigma} \cdot \boldsymbol{L} & 0 \\
0 & \boldsymbol{\sigma} \cdot \boldsymbol{L}
\end{array}\right] \\
& =\left[\begin{array}{cc}
0 & (\boldsymbol{\sigma} \cdot \boldsymbol{p})(\boldsymbol{\sigma} \cdot \boldsymbol{L}) \\
(\boldsymbol{\sigma} \cdot \boldsymbol{p})(\boldsymbol{\sigma} \cdot \boldsymbol{L}) & 0
\end{array}\right]  \tag{3.121}\\
(\boldsymbol{\sigma} \cdot \boldsymbol{L})(\boldsymbol{\alpha} \cdot \boldsymbol{p}) & =\left[\begin{array}{cc}
\boldsymbol{\sigma} \cdot \boldsymbol{L} & 0 \\
0 & \boldsymbol{\sigma} \cdot \boldsymbol{L}
\end{array}\right]\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \boldsymbol{p} \\
\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0
\end{array}\right] \\
& =\left[\begin{array}{cc}
0 & (\boldsymbol{\sigma} \cdot \boldsymbol{L})(\boldsymbol{\sigma} \cdot \boldsymbol{p}) \\
(\boldsymbol{\sigma} \cdot \boldsymbol{L})(\boldsymbol{\sigma} \cdot \boldsymbol{p}) & 0
\end{array}\right] \tag{3.122}
\end{align*}
$$

Let us evaluate $(\boldsymbol{\sigma} \cdot \boldsymbol{p})(\boldsymbol{\sigma} \cdot \boldsymbol{L})$ and $(\boldsymbol{\sigma} \cdot \boldsymbol{L})(\boldsymbol{\sigma} \cdot \boldsymbol{p})$.

$$
\begin{align*}
& (\boldsymbol{\sigma} \cdot \boldsymbol{p})(\boldsymbol{\sigma} \cdot \boldsymbol{L})=\boldsymbol{p} \cdot \boldsymbol{L}+i \boldsymbol{\sigma} \cdot(\boldsymbol{p} \times \boldsymbol{L})  \tag{3.123}\\
& (\boldsymbol{\sigma} \cdot \boldsymbol{L})(\boldsymbol{\sigma} \cdot \boldsymbol{p})=\boldsymbol{L} \cdot \boldsymbol{p}+i \boldsymbol{\sigma} \cdot(\boldsymbol{L} \times \boldsymbol{p}) \tag{3.124}
\end{align*}
$$

Since $\boldsymbol{L}=\boldsymbol{r} \times \boldsymbol{p}$, the vector $\boldsymbol{L}$ is perpendicular to the plane containing $\boldsymbol{r}$ and $\boldsymbol{p}$. Hence

$$
\boldsymbol{p} \cdot \boldsymbol{L}=\boldsymbol{L} \cdot \boldsymbol{p}=0
$$

The quantities $\boldsymbol{\sigma} \cdot(\boldsymbol{p} \times \boldsymbol{L})$ and $\boldsymbol{\sigma} \cdot(\boldsymbol{L} \times \boldsymbol{p})$ can be evaluated using the commutation relations between the components of $\boldsymbol{L}$ and components of $\boldsymbol{p}$. In natural units $(\hbar=c=1)$,

$$
\begin{align*}
& {\left[L_{x}, p_{x}\right]_{-}=\left[L_{y}, p_{y}\right]_{-}=\left[L_{z}, p_{z}\right]_{-}=0}  \tag{3.125}\\
& {\left[L_{x}, p_{y}\right]_{-}=i p_{z}, \quad\left[L_{y}, p_{z}\right]_{-}=i p_{x}, \quad\left[L_{z}, p_{x}\right]_{-}=i p_{y}} \tag{3.126}
\end{align*}
$$

$\boldsymbol{\sigma} \cdot(\boldsymbol{p} \times \boldsymbol{L})=\sigma_{x}\left(p_{y} L_{z}-p_{z} L_{y}\right)+\sigma_{y}\left(p_{z} L_{x}-p_{x} L_{z}\right)+\sigma_{z}\left(p_{x} L_{y}-p_{y} L_{x}\right)$ $\boldsymbol{\sigma} \cdot(\boldsymbol{L} \times \boldsymbol{p})=\sigma_{x}\left(L_{y} p_{z}-L_{z} p_{y}\right)+\sigma_{y}\left(L_{z} p_{x}-L_{x} p_{z}\right)+\sigma_{z}\left(L_{x} p_{y}-L_{y} p_{x}\right)$.

Adding, we get

$$
\begin{aligned}
\boldsymbol{\sigma} \cdot(\boldsymbol{p} \times \boldsymbol{L})+\boldsymbol{\sigma} \cdot(\boldsymbol{L} \times \boldsymbol{p})= & \sigma_{x}\left\{\left[L_{y}, p_{z}\right]_{-}-\left[L_{z}, p_{y}\right]_{-}\right\} \\
& +\sigma_{y}\left\{\left[L_{z}, p_{x}\right]_{-}-\left[L_{x}, p_{z}\right]_{-}\right\} \\
& +\sigma_{z}\left\{\left[L_{x}, p_{y}\right]_{-}-\left[L_{y}, p_{x}\right]_{-}\right\} \\
= & \sigma_{x}\left\{i p_{x}+i p_{x}\right\}+\sigma_{y}\left\{i p_{y}+i p_{y}\right\} \\
& +\sigma_{z}\left\{i p_{z}+i p_{z}\right\} \\
= & 2 i \boldsymbol{\sigma} \cdot \boldsymbol{p} .
\end{aligned}
$$

Collecting he above results and substituting in Eq. (3.120), we get

$$
\begin{align*}
{[\boldsymbol{\alpha} \cdot \boldsymbol{p}, \beta \boldsymbol{\sigma} \cdot \boldsymbol{L}]_{-} } & =-\beta\{(\boldsymbol{\alpha} \cdot \boldsymbol{p})(\boldsymbol{\sigma} \cdot \boldsymbol{L})+(\boldsymbol{\sigma} \cdot \boldsymbol{L})(\boldsymbol{\alpha} \cdot \boldsymbol{p})\} \\
& =-\left[\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right]\left[\begin{array}{cc}
0 & -2 \hbar \boldsymbol{\sigma} \cdot \boldsymbol{p} \\
-2 \hbar \boldsymbol{\sigma} \cdot \boldsymbol{p} & 0
\end{array}\right] \\
& =2 \hbar \boldsymbol{\sigma} \cdot \boldsymbol{p}\left[\begin{array}{rr}
0 & 1 \\
-1 & 0
\end{array}\right]  \tag{3.127}\\
{[\boldsymbol{\alpha} \cdot \boldsymbol{p}, \beta]_{-} } & =\boldsymbol{\alpha} \cdot \boldsymbol{p} \beta-\beta \boldsymbol{\alpha} \cdot \boldsymbol{p} \\
& =-\beta \boldsymbol{\alpha} \cdot \boldsymbol{p}-\beta \boldsymbol{\alpha} \cdot \boldsymbol{p} \\
& =-2 \beta \boldsymbol{\alpha} \cdot \boldsymbol{p} \\
& =-2\left[\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right]\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \boldsymbol{p} \\
\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0
\end{array}\right] \\
& =-2\left[\begin{array}{rr}
0 & \boldsymbol{\sigma} \cdot \boldsymbol{p} \\
-\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0
\end{array}\right] . \tag{3.128}
\end{align*}
$$

Equations (3.127) and (3.128) together yield zero for the commutator $[\boldsymbol{\alpha} \cdot \boldsymbol{p}, K]_{-}$.

## To show that $[J, K]_{-}$

Since $\boldsymbol{J}=\boldsymbol{L}+\frac{1}{2} \boldsymbol{\sigma}$, we need to evaluate the commutator

$$
\left[\boldsymbol{L}+\frac{1}{2} \boldsymbol{\sigma}, \beta(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1)\right]_{-}
$$

$\beta$ commutes with both $\boldsymbol{L}$ and $\boldsymbol{\sigma}$. So, we need to evaluate only the two commutators

$$
[\boldsymbol{L}, \beta \boldsymbol{\sigma} \cdot \boldsymbol{L}]_{-} \quad \text { and } \quad\left[\frac{1}{2} \boldsymbol{\sigma}, \beta \boldsymbol{\sigma} \cdot \boldsymbol{L}\right]_{-}
$$

Let us first evaluate the commutator $[\boldsymbol{L}, \beta \boldsymbol{\sigma} \cdot \boldsymbol{L}]_{-}$.

$$
\begin{align*}
{[L, \beta \boldsymbol{\sigma} \cdot \boldsymbol{L}]_{-} } & =\boldsymbol{L} \beta \boldsymbol{\sigma} \cdot \boldsymbol{L}-\beta \boldsymbol{\sigma} \cdot \boldsymbol{L} \\
& =\beta\{\boldsymbol{L}(\boldsymbol{\sigma} \cdot \boldsymbol{L})-(\boldsymbol{\sigma} \cdot \boldsymbol{L}) \boldsymbol{L}\} \tag{3.129}
\end{align*}
$$

Let us expand $\boldsymbol{L}(\boldsymbol{\sigma} \cdot \boldsymbol{L})-(\boldsymbol{\sigma} \cdot \boldsymbol{L}) \boldsymbol{L}$ in terms of their components.

$$
\begin{align*}
& \boldsymbol{L}(\boldsymbol{\sigma} \cdot \boldsymbol{L})-(\boldsymbol{\sigma} \cdot \boldsymbol{L}) \boldsymbol{L} \\
&=\left(L_{x} \boldsymbol{e}_{x}+L_{y} \boldsymbol{e}_{y}+L_{z} \boldsymbol{e}_{z}\right)\left(\sigma_{x} L_{x}+\sigma_{y} L_{y}+\sigma_{z} L_{z}\right) \\
&-\left(\sigma_{x} L_{x}+\sigma_{y} L_{y}+\sigma_{z} L_{z}\right)\left(L_{x} \boldsymbol{e}_{x}+L_{y} \boldsymbol{e}_{y}+L_{z} \boldsymbol{e}_{z}\right) \\
&=\left(\sigma_{x} L_{x}^{2}+\sigma_{y} L_{x} L_{y}+\sigma_{z} L_{x} L_{z}-\sigma_{x} L_{x}^{2}-\sigma_{y} L_{y} L_{x}-\sigma_{z} L_{z} L_{x}\right) \boldsymbol{e}_{x} \\
&+(\cdots) \boldsymbol{e}_{y}+(\cdots) \boldsymbol{e}_{z} \\
&=\left\{\sigma_{y}\left[L_{x}, L_{y}\right]_{-}+\sigma_{z}\left[L_{x}, L_{z}\right]_{-}\right\} \boldsymbol{e}_{x}+(\cdots) \boldsymbol{e}_{y}+(\cdots) \boldsymbol{e}_{z} \\
&= i\left(\sigma_{y} L_{z}-\sigma_{z} L_{y}\right) \boldsymbol{e}_{x}+(\cdots) \boldsymbol{e}_{y}+(\cdots) \boldsymbol{e}_{z} \\
&= i(\boldsymbol{\sigma} \times \boldsymbol{L}) \tag{3.130}
\end{align*}
$$

Thereby, we arrive at a compact expression for the commutator.

$$
\begin{equation*}
[\boldsymbol{L}, \beta \boldsymbol{\sigma} \cdot \boldsymbol{L}]_{-}=i \beta(\boldsymbol{\sigma} \times \boldsymbol{L}) \tag{3.131}
\end{equation*}
$$

We shall now evaluate the other commutator .

$$
\begin{align*}
{[\boldsymbol{\sigma}, \beta \boldsymbol{\sigma} \cdot \boldsymbol{L}]_{-} } & =\boldsymbol{\sigma} \beta(\boldsymbol{\sigma} \cdot \boldsymbol{L})-\beta(\boldsymbol{\sigma} \cdot \boldsymbol{L}) \boldsymbol{\sigma} \\
& =\beta\{\boldsymbol{\sigma}(\boldsymbol{\sigma} \cdot \boldsymbol{L})-(\boldsymbol{\sigma} \cdot \boldsymbol{L}) \boldsymbol{\sigma}\} \tag{3.132}
\end{align*}
$$

Using the relations

$$
\begin{align*}
\boldsymbol{\sigma}(\boldsymbol{\sigma} \cdot \boldsymbol{L}) & =\boldsymbol{L}-i \boldsymbol{\sigma} \times \boldsymbol{L}  \tag{3.133}\\
(\boldsymbol{\sigma} \cdot \boldsymbol{L}) \boldsymbol{\sigma} & =\boldsymbol{L}+i \boldsymbol{\sigma} \times \boldsymbol{L} \tag{3.134}
\end{align*}
$$

we find

$$
[\boldsymbol{\sigma}, \beta \boldsymbol{\sigma} \cdot \boldsymbol{L}]_{-}=-2 i \beta(\boldsymbol{\sigma} \times \boldsymbol{L}) ; \quad\left[\frac{1}{2} \boldsymbol{\sigma}, \beta \boldsymbol{\sigma} \times \boldsymbol{L}\right]_{-}=-i \beta(\boldsymbol{\sigma} \times \boldsymbol{L})(3.135)
$$

Thus the contribution from the commutator (3.131) is exactly canceled by the commutator (3.135). Thereby we are assured that the operators $\boldsymbol{J}$ and $K$ commute, permitting us to find the simultaneous eigenfunctions for both the operators.
3.3 It is simpler to prove the inverse relation.

$$
\begin{aligned}
&-i(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \boldsymbol{L})=\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}} \times \boldsymbol{L} \\
&-i(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \boldsymbol{L})= \\
&-i\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} \\
\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} & 0
\end{array}\right]\left[\begin{array}{cc}
\boldsymbol{\sigma} \cdot \boldsymbol{L} & 0 \\
0 & \boldsymbol{\sigma} \cdot \boldsymbol{L}
\end{array}\right] \\
&=-i\left[\begin{array}{cc}
0 & (\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \boldsymbol{L}) \\
(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \boldsymbol{L}) & 0
\end{array}\right] \\
&=-i\left[\begin{array}{cc}
0 & i \boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} \times \boldsymbol{L} \\
i \boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} \times \boldsymbol{L} & 0
\end{array}\right] \\
&=\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}} \times \boldsymbol{L} .
\end{aligned}
$$

The above result is obtained by using the relation

$$
\begin{aligned}
(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \boldsymbol{L}) & =\hat{\boldsymbol{r}} \cdot \boldsymbol{L}+i \boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} \times \boldsymbol{L} \\
& =\hat{\boldsymbol{r}} \cdot \boldsymbol{r} \times \boldsymbol{p}+i \boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} \times \boldsymbol{L} \\
& =i \boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} \times \boldsymbol{L},
\end{aligned}
$$

since $\hat{\boldsymbol{r}} \cdot \boldsymbol{r} \times \boldsymbol{p}=\hat{\boldsymbol{r}} \times \boldsymbol{r} \cdot \boldsymbol{p}=0$.
3.4 Let us consider the operator $(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})^{2}$.

$$
\begin{aligned}
(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})^{2} & =(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})(\alpha \cdot \hat{\boldsymbol{r}}) \\
& =\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} \\
\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} & 0
\end{array}\right]\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} \\
\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} & 0
\end{array}\right] \\
& =\left[\begin{array}{cc}
(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}}) & 0 \\
0 & (\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})
\end{array}\right]
\end{aligned}
$$

Since

$$
(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})=1
$$

we find that $(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})^{2}$ is a unit matrix. Hence $\boldsymbol{\alpha} \cdot \boldsymbol{p}$ can be written as $(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})^{2}(\boldsymbol{\alpha} \cdot \boldsymbol{p})$. First let us find

$$
\begin{aligned}
(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\alpha} \cdot \boldsymbol{p}) & =\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} \\
\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}} & 0
\end{array}\right]\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \cdot \boldsymbol{p} \\
\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0
\end{array}\right] \\
& =\left[\begin{array}{cc}
(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \boldsymbol{p}) & 0 \\
0 & (\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \boldsymbol{p})
\end{array}\right]
\end{aligned}
$$

We know that

$$
(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \boldsymbol{p})=\hat{\boldsymbol{r}} \cdot \boldsymbol{p}+i \boldsymbol{\sigma} \cdot(\hat{\boldsymbol{r}} \times \boldsymbol{p})=\hat{\boldsymbol{r}} \cdot \boldsymbol{p}+\frac{i}{r}(\boldsymbol{\sigma} \cdot \boldsymbol{L})
$$

Replacing $\boldsymbol{\sigma} \cdot \boldsymbol{L}$ by the operator $\beta K-1$, we find

$$
(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \boldsymbol{p})=\hat{\boldsymbol{r}} \cdot \boldsymbol{p}+\frac{i}{r}(\beta K-1)
$$

Since $(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})^{2}=1$, we get

$$
\begin{aligned}
\boldsymbol{\alpha} \cdot \boldsymbol{p} & =(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})^{2}(\boldsymbol{\alpha} \cdot \boldsymbol{p}) \\
& =(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})\left\{\hat{\boldsymbol{r}} \cdot \boldsymbol{p}+\frac{i}{r}(\beta K-1)\right\} \\
& =-i(\boldsymbol{\alpha} \cdot \hat{\boldsymbol{r}})\left\{\frac{\partial}{\partial r}-\frac{1}{r}(\beta K-1)\right\}
\end{aligned}
$$

since $\hat{\boldsymbol{r}} \cdot \boldsymbol{p}=-i \hat{\boldsymbol{r}} \cdot \boldsymbol{\nabla}=-i \frac{\partial}{\partial r}$.
3.5 Let us consider an operator $T_{k}(1) \cdot T_{k}(2)$, which is a scalar product of two tensors, each of rank $k$. The matrix element of this operator, taken between two coupled angular momentum states $\left|j_{1}, j_{2}, j, m\right\rangle$ and $\left|j_{1}^{\prime}, j_{2}^{\prime}, j^{\prime}, m^{\prime}\right\rangle$ is given by

$$
\begin{align*}
Q & =\left\langle j_{1}^{\prime} j_{2}^{\prime} j^{\prime} m^{\prime}\right| T_{k}(1) \cdot T_{k}(2)\left|j_{1} j_{2} j m\right\rangle \\
& =\sum_{\mu}(-1)^{\mu}\left\langle j_{1}^{\prime} j_{2}^{\prime} j m\right| T_{k}^{\mu}(1) T_{k}^{-\mu}(2)\left|j_{1} j_{2} j m\right\rangle \tag{3.136}
\end{align*}
$$

Since the operator is a scalar in the coupled angular momentum state, $j=j^{\prime}$ and $m=m^{\prime}$.
Using the angular momentum algebra ${ }^{8}$, the matrix element can be written as a product of two uncoupled angular momentum matrix elements, using the U-coefficient or Racah coefficient.

$$
\begin{equation*}
Q=(-1)^{k} \frac{\left[j_{2}^{\prime}\right]}{\left[j_{2}\right]} U\left(j_{1} k j j_{2}^{\prime}, j_{1}^{\prime} j_{2}\right)\left\langle j_{1}^{\prime}\left\|T_{k}(1)\right\| j_{1}\right\rangle\left\langle j_{2}^{\prime}\left\|T_{k}(2)\right\| j_{2}\right\rangle \tag{3.137}
\end{equation*}
$$

[^21]where the U-coefficient is expressed in terms of the Racah coefficient
$$
U(a b c d, e f)=[e][f] W(a b c d, e f)
$$
with the symbol [ $j$ ] standing for $(2 j+1)^{1 / 2}$.
With this brief introduction to the angular momentum algebra, let us consider the evaluation of the matrix element
$$
\left\langle l_{f} \frac{1}{2} j_{f} m_{f}\right| \boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}}\left|l_{i} \frac{1}{2} j_{i} m_{i}\right\rangle
$$

Since $\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}}$ is a scalar (strictly a pseudoscalar) in the combined space of configuration and spin, $j_{f}$ should be equal to $j_{i}$ and $m_{i}=m_{f}$. So, let us impose the condition $j_{i}=j_{f}=j$. Since the operator is a scalar, there will be no dependence on the magnetic quantum number.

$$
\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}}=\sqrt{\frac{4 \pi}{3}} \boldsymbol{\sigma} \cdot Y_{1}(\hat{\boldsymbol{r}})
$$

Applying the result (3.137), we obtain

$$
\begin{aligned}
\left\langle l_{f} \frac{1}{2} j\|\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}}\| l_{i} \frac{1}{2} j\right\rangle & =-U\left(l_{i} 1 j \frac{1}{2}, l_{f} \frac{1}{2}\right)\left\langle l_{f}\left\|Y_{1}\right\| l_{i}\right\rangle\left\langle\frac{1}{2}\|\boldsymbol{\sigma}\| \frac{1}{2}\right\rangle \\
& =-\left[l_{f}\right]\left[\frac{1}{2}\right] W\left(l_{i} 1 j \frac{1}{2}, l_{f} \frac{1}{2}\right)\left\langle l_{f}\left\|Y_{1}\right\| l_{i}\right\rangle\left\langle\frac{1}{2}\|\boldsymbol{\sigma}\| \frac{1}{2}\right\rangle,(3.138)
\end{aligned}
$$

with

$$
\begin{aligned}
\left\langle l_{f}\left\|Y_{1}\right\| l_{i}\right\rangle & =\frac{\left[l_{i}\right][1]}{\sqrt{4 \pi}\left[l_{f}\right]}\left[\begin{array}{ccc}
l_{i} & 1 & l_{f} \\
0 & 0 & 0
\end{array}\right] \\
\left\langle\frac{1}{2}\|\boldsymbol{\sigma}\| \frac{1}{2}\right\rangle & =[1]
\end{aligned}
$$

Let us substitute the algebraic expressions for the W-coefficient and the C.G. coefficient which are available in Tabular form ${ }^{9}$.

Case 1: $l_{i}=j-\frac{1}{2}, \quad l_{f}=j+\frac{1}{2}$

$$
\begin{aligned}
& W\left(l_{i} j 1 \frac{1}{2}, \frac{1}{2} l_{f}\right)=\left\{\frac{1}{3(2 j+1)}\right\}^{\frac{1}{2}} \\
& {\left[\begin{array}{ccc}
l_{i} & 1 & l_{f} \\
0 & 0 & 0
\end{array}\right]=\left\{\frac{j+\frac{1}{2}}{2 j}\right\}^{\frac{1}{2}}}
\end{aligned}
$$

Substituting these values in Eq. (3.138), we get

$$
\left\langle l_{f} \frac{1}{2} j\|\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}}\| l_{i} \frac{1}{2} j\right\rangle=-1
$$

[^22]Case 2: $l_{i}=j+\frac{1}{2}, l_{f}=j-\frac{1}{2}$

$$
\begin{aligned}
& W\left(l_{i} j 1 \frac{1}{2}, \frac{1}{2} l_{f}\right)=-\left\{\frac{1}{3(2 j+1)}\right\}^{\frac{1}{2}} ; \\
& {\left[\begin{array}{ccc}
l_{i} & 1 & l_{f} \\
0 & 0 & 0
\end{array}\right]=-\left\{\frac{j+\frac{1}{2}}{2 j+2}\right\}^{\frac{1}{2}} .}
\end{aligned}
$$

Substituting these values in Eq. (3.138), we get

$$
\left\langle l_{f} \frac{1}{2} j\|\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}}\| l_{i} \frac{1}{2} j\right\rangle=-1 .
$$

Then, it follows that

$$
\begin{aligned}
\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}}\left|l_{a} \frac{1}{2} j m\right\rangle & =-\left|l_{b} \frac{1}{2} j m\right\rangle \\
\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}}\left|l_{b} \frac{1}{2} j m\right\rangle & =-\left|l_{a} \frac{1}{2} j m\right\rangle .
\end{aligned}
$$

The above result can be obtained from a simple consideration. Since $\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}}$ is a pseudo-scalar, its parity is -1 . Therefore

$$
l_{f}=l_{i} \pm 1
$$

In other words, if $l_{i}=j+\frac{1}{2}$, then $l_{f}=j-\frac{1}{2}$ or vice versa.
Since the square of the operator

$$
(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})(\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}})=\hat{\boldsymbol{r}} \cdot \hat{\boldsymbol{r}}+i \boldsymbol{\sigma} \cdot(\hat{\boldsymbol{r}} \times \hat{\boldsymbol{r}})=1
$$

it follows that

$$
\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}}=-1
$$

remembering that $\boldsymbol{\sigma} \cdot \hat{\boldsymbol{r}}$ is a pseudo-scalar.
3.6 (a) In Classical Mechanics,

$$
\begin{aligned}
\boldsymbol{L} & =\boldsymbol{r} \times \boldsymbol{p} \\
\boldsymbol{L}^{2} & =(\boldsymbol{r} \times \boldsymbol{p}) \cdot(\boldsymbol{r} \times \boldsymbol{p})
\end{aligned}
$$

Using the vector algebra,

$$
(\boldsymbol{A} \times \boldsymbol{B}) \cdot(\boldsymbol{C} \times \boldsymbol{D})=(\boldsymbol{A} \cdot \boldsymbol{C})(\boldsymbol{B} \cdot \boldsymbol{D})-(\boldsymbol{A} \cdot \boldsymbol{D})(\boldsymbol{B} \cdot \boldsymbol{C})
$$

we get

$$
\begin{equation*}
\boldsymbol{L}^{2}=\boldsymbol{r}^{2} \boldsymbol{p}^{2}-(\boldsymbol{r} \cdot \boldsymbol{p})^{2} \tag{3.139}
\end{equation*}
$$

Thus, we get

$$
\begin{equation*}
\boldsymbol{p}^{2}=\frac{1}{r^{2}} \boldsymbol{L}^{2}+p_{r}^{2} \tag{3.140}
\end{equation*}
$$

if $p_{r} \equiv \frac{1}{r}(\boldsymbol{r} \cdot \boldsymbol{p})$ is defined in classical mechanics as the radial component of momentum in central field.

In Quantum Mechanics, $\boldsymbol{r}$ and $\boldsymbol{p}$ do not commute. Using the commutation relations between the conjugate variables,

$$
\left[x, p_{x}\right]_{-}=\left[y, p_{y}\right]_{-}=\left[z, p_{z}\right]_{-}=i \hbar
$$

we get an extra term $i \hbar(\boldsymbol{r} \cdot \boldsymbol{p})$.

$$
\begin{equation*}
\boldsymbol{L}^{2}=\boldsymbol{r}^{2} \boldsymbol{p}^{2}-(\boldsymbol{r} \cdot \boldsymbol{p})^{2}+i \hbar(\boldsymbol{r} \cdot \boldsymbol{p}) \tag{3.141}
\end{equation*}
$$

## To derive Eq. (3.141)

Since $\boldsymbol{L}^{2}=(\boldsymbol{r} \times \boldsymbol{p}) \cdot(\boldsymbol{r} \times \boldsymbol{p})$, we get

$$
\begin{align*}
L_{x}^{2} & =\left(y p_{z}-z p_{y}\right)\left(y p_{z}-z p_{y}\right) \\
& =y^{2} p_{z}^{2}+z^{2} p_{y}^{2}-2 y p_{y} z p_{z}+i \hbar\left(y p_{y}+z p_{z}\right) \tag{3.142}
\end{align*}
$$

using the commutation relations between the conjugate variables. Similarly we can find the other components.

$$
\begin{align*}
L_{y}^{2} & =z^{2} p_{x}^{2}+x^{2} p_{z}^{2}-2 z p_{x} x p_{x}+i \hbar\left(z p_{z}+x p_{x}\right)  \tag{3.143}\\
L_{z}^{2} & =x^{2} p_{y}^{2}+y^{2} p_{x}^{2}-2 x p_{x} y p_{y}+i \hbar\left(x p_{x}+y p_{y}\right) \tag{3.144}
\end{align*}
$$

Adding, we get

$$
\begin{align*}
\boldsymbol{L}^{2}= & L_{x}^{2}+L_{y}^{2}+L_{z}^{2} \\
= & x^{2}\left(p_{x}^{2}+p_{y}^{2}+p_{z}^{2}\right)-x^{2} p_{x}^{2}+y^{2}\left(p_{x}^{2}+p_{y}^{2}+p_{z}^{2}\right)-y^{2} p_{y}^{2} \\
& +z^{2}\left(p_{x}^{2}+p_{y}^{2}+p_{z}^{2}\right)-z^{2} p_{z}^{2}-2 y p_{y} z p_{z}-2 z p_{z} x p_{x}-2 x p_{x} y p_{y} \\
& +2 i \hbar\left(x p_{x}+y p_{y}+z p_{z}\right) \\
= & \boldsymbol{r}^{2} \boldsymbol{p}^{2}-\left(x^{2} p_{x}^{2}+y^{2} p_{y}^{2}+z^{2} p_{z}^{2}\right)-2 x p_{x} y p_{y}-2 y p_{y} z p_{z}-2 z p_{z} x p_{x} \\
& +2 i \hbar(\boldsymbol{r} \cdot \boldsymbol{p}) . \tag{3.145}
\end{align*}
$$

In a similar way, we can find

$$
\begin{align*}
(\boldsymbol{r} \cdot \boldsymbol{p})^{2}= & \left(x p_{x}+y p_{y}+z p_{z}\right)\left(x p_{x}+y p_{y}+z p_{z}\right) \\
= & x^{2} p_{x}^{2}+y^{2} p_{y}^{2}+z^{2} p_{z}^{2}-i \hbar(\boldsymbol{r} \cdot \boldsymbol{p}) \\
& +2 x p_{x} y p_{y}+2 y p_{y} z p_{z}+2 z p_{z} x p_{x} \tag{3.146}
\end{align*}
$$

Using (3.146) in Eq. (3.145), we get Eq. (3.141).

$$
\boldsymbol{L}^{2}=\boldsymbol{r}^{2} \boldsymbol{p}^{2}-(\boldsymbol{r} \cdot \boldsymbol{p})^{2}+i \hbar(\boldsymbol{r} \cdot \boldsymbol{p})
$$

## To deduce an expression for $p_{r}$ in Quantum Mechanics

The symbol $p_{r}$ denotes the radial component of momentum in a central potential (spherically symmetric potential). So, it is convenient to work in spherical coordinates. From Eq. (3.141), we obtain

$$
\begin{align*}
\boldsymbol{p}^{2} & =\frac{1}{r^{2}}\left\{\boldsymbol{L}^{2}+(\boldsymbol{r} \cdot \boldsymbol{p})^{2}-i \hbar \boldsymbol{r} \cdot \boldsymbol{p}\right\} \\
& =\frac{1}{r^{2}}\left\{\boldsymbol{L}^{2}-\hbar^{2}\left(r \frac{\partial}{\partial r}\right)^{2}-\hbar^{2} r \frac{\partial}{\partial r}\right\} \\
& =\frac{1}{r^{2}} \boldsymbol{L}^{2}-\hbar^{2}\left\{\frac{\partial^{2}}{\partial r^{2}}+\frac{2}{r} \frac{\partial}{\partial r}\right\} \tag{3.147}
\end{align*}
$$

since

$$
\boldsymbol{r} \cdot \boldsymbol{p}=-i \hbar \boldsymbol{r} \cdot \boldsymbol{\nabla}=-i \hbar r \frac{\partial}{\partial r}
$$

and

$$
\left(r \frac{\partial}{\partial r}\right)^{2}=r \frac{\partial}{\partial r}\left(r \frac{\partial}{\partial r}\right)=r \frac{\partial}{\partial r}+r^{2} \frac{\partial^{2}}{\partial r^{2}}
$$

Equation (3.147) can be written in a form similar to (3.140)

$$
\begin{equation*}
\boldsymbol{p}^{2}=\frac{1}{r^{2}} \boldsymbol{L}^{2}+p_{r}^{2} \tag{3.148}
\end{equation*}
$$

if $p_{r}$ in quantum mechanics is defined as

$$
\begin{equation*}
p_{r}=\frac{1}{r}(\boldsymbol{r} \cdot \boldsymbol{p}-i \hbar) \tag{3.149}
\end{equation*}
$$

Writing $\boldsymbol{p}$ as a differential operator, $\boldsymbol{r} \cdot \boldsymbol{p}=-i \hbar \boldsymbol{r} \cdot \boldsymbol{\nabla}=-i \hbar r \frac{\partial}{\partial r}$. It can be easily verified that

$$
\begin{equation*}
p_{r}=-i \hbar\left(\frac{\partial}{\partial r}+\frac{1}{r}\right) \tag{3.150}
\end{equation*}
$$

and

$$
\begin{align*}
p_{r}^{2} & =-\hbar^{2}\left(\frac{\partial}{\partial r}+\frac{1}{r}\right)\left(\frac{\partial}{\partial r}+\frac{1}{r}\right) \\
& =-\hbar^{2}\left\{\frac{\partial^{2}}{\partial r^{2}}+\frac{2}{r} \frac{\partial}{\partial r}\right\} \tag{3.151}
\end{align*}
$$

Substituting (3.151) into Eq. (3.148), we get back Eq. (3.147).
(b) It can be easily verified that $r$ and $p_{r}$ obey the commutation relation

$$
\left[r, p_{r}\right]_{-}=i \hbar
$$

using Eq. (3.150) for $p_{r}$.
(c) It can be shown that $p_{r}$ as given by Eq. (3.150) is a Hermitian operator. If $\Phi$ and $\Psi$ are two arbitrary state vectors, the the condition for Hermiticity of $p_{r}$ is

$$
\begin{equation*}
\int_{0}^{\infty} d r r^{2} \Phi^{*} p_{r} \Psi=\int_{0}^{\infty} d r r^{2}\left(p_{r} \Phi\right)^{*} \Psi \tag{3.152}
\end{equation*}
$$

Let us consider separately the two terms occurring in Eq. (3.150) for $p_{r}$. First term:

$$
\begin{equation*}
\int d r r^{2} \Phi^{*}\left(-i \hbar \frac{\partial}{\partial r}\right) \Psi=-\left.i \hbar r^{2} \Phi^{*} \Psi\right|_{0} ^{\infty}-\int d r \Psi\left(-i \hbar \frac{\partial}{\partial r} r^{2} \Phi^{*}\right) \tag{3.153}
\end{equation*}
$$

In Eq. (3.153), the first term vanishes when the limits are taken and the second term alone survives and it is equal to

$$
\begin{equation*}
-\int d r r^{2}\left(i \hbar \frac{\partial}{\partial r} \Phi\right)^{*} \Psi-2 \int d r r^{2}\left(i \hbar \frac{1}{r} \Phi\right)^{*} \Psi \tag{3.154}
\end{equation*}
$$

Second term:

$$
\begin{equation*}
\int d r r^{2} \Phi^{*}\left(-i \hbar \frac{1}{r}\right) \Psi=\int d r r^{2}\left(i \hbar \frac{1}{r} \Phi\right)^{*} \Psi \tag{3.155}
\end{equation*}
$$

Adding the contributions from Eqs. (3.153) and (3.155), we establish the Hermiticity condition (3.152).
3.7 Given the operators

$$
\alpha_{r}=\frac{1}{r}(\boldsymbol{\alpha} \cdot \boldsymbol{r}) ; \quad p_{r}=\frac{1}{r}(\boldsymbol{r} \cdot \boldsymbol{p}-i) ; \quad K=\beta(\boldsymbol{\sigma} \cdot \boldsymbol{L}+1)
$$

let us express $(\boldsymbol{\alpha} \cdot \boldsymbol{r})(\boldsymbol{\alpha} \cdot \boldsymbol{p})$ in terms of them.

$$
\begin{array}{rlrl}
(\boldsymbol{\alpha} \cdot \boldsymbol{r})(\boldsymbol{\alpha} \cdot \boldsymbol{p}) & =\boldsymbol{r} \cdot \boldsymbol{p}+i \boldsymbol{\sigma} \cdot(\boldsymbol{r} \times \boldsymbol{p}) & \\
& =\left(r p_{r}+i\right)+i \boldsymbol{\sigma} \cdot \boldsymbol{L}, & & \text { since } p_{r}=\frac{1}{r}(\boldsymbol{r} \cdot \boldsymbol{p}-i) \\
& =\left(r p_{r}+i\right)+i(\beta K-1), & & \text { since } \beta^{2}=1 \\
& =r p_{r}+i \beta K . & &
\end{array}
$$

Alternatively, $(\boldsymbol{\alpha} \cdot \boldsymbol{r})(\boldsymbol{\alpha} \cdot \boldsymbol{p})$ can be written as

$$
(\boldsymbol{\alpha} \cdot \boldsymbol{r})(\boldsymbol{\alpha} \cdot \boldsymbol{p})=r \alpha_{r}(\boldsymbol{\alpha} \cdot \boldsymbol{p})
$$

So,

$$
r \alpha_{r}(\boldsymbol{\alpha} \cdot \boldsymbol{p})=r p_{r}+i \beta K
$$

Multiplying both sides by $\alpha_{r}$, we get

$$
\boldsymbol{\alpha} \cdot \boldsymbol{p}=\alpha_{r} p_{r}+\frac{i}{r} \alpha_{r} \beta K
$$

since

$$
\alpha_{r}^{2}=\frac{1}{r^{2}}(\boldsymbol{\alpha} \cdot \boldsymbol{r})(\boldsymbol{\alpha} \cdot \boldsymbol{r})=\frac{1}{r^{2}}\{\boldsymbol{r} \cdot \boldsymbol{r}+i \boldsymbol{\sigma} \cdot(\boldsymbol{r} \times \boldsymbol{r})\}=1 .
$$

Using the above results, the Dirac Hamiltonian can be written as

$$
H=\alpha_{r} p_{r}+\frac{i}{r} \alpha_{r} \beta K+\beta m+V(r)
$$

## Chapter 4

## The Neutrino

The neutrino was first postulated by W. Pauli in 1930 in order to explain the energy, momentum and angular momentum conservation in beta decay. It is a charge-less, massless particle with spin $-\frac{1}{2}$ which could not be detected in the early $\beta$ - decay experiments but conjectured to participate in the reaction in order to preserve the conservation laws which are held sacrosanct and form the main basis for all physical sciences. Pauli was very apologetic for making this postulate and is supposed to have stated "I have done a terrible thing. I have postulated a particle that cannot be detected." However, after 25 years, Frederic Reines and Clyde Lorrain Cowan,Jr. ${ }^{1}$ succeeded, in the year 1956, in devising and performing an experiment by which the neutrino was detected. Reines shared the Nobel Prize in Physics awarded in 1995 along with M. L. Perl ${ }^{2}$ who discovered the Tau lepton in 1975 at the Stanford Linear Accelerator Laboratory.

### 4.1 The Dirac equation for the neutrino

Since neutrino is a spin $-\frac{1}{2}$ particle, it should obey the free particle Dirac equation

$$
\begin{equation*}
\left(c \boldsymbol{\alpha} \cdot \boldsymbol{p}+\beta m c^{2}\right) \psi=E \psi, \tag{4.1}
\end{equation*}
$$

[^23]for which the solutions are given by Eqs. (2.16) and (2.17) in chapter 2.
\[

$$
\begin{array}{ll}
\psi_{1}=N\left[\begin{array}{c}
1 \\
0 \\
\frac{p_{z}}{W+\mu} \\
\frac{p_{+}}{W+\mu}
\end{array}\right] ; & \psi_{2}=N\left[\begin{array}{c}
0 \\
1 \\
\frac{p_{-}}{W+\mu} \\
\frac{-p_{z}}{W+\mu}
\end{array}\right] ; \\
\psi_{3}=N\left[\begin{array}{c}
\frac{-p_{z}}{W+\mu} \\
\frac{-p_{+}}{W+\mu} \\
1 \\
0
\end{array}\right] ; & \psi_{4}=N\left[\begin{array}{c}
\frac{-p_{-}}{W+\mu} \\
\frac{p_{z}}{W+\mu} \\
0 \\
1
\end{array}\right] ; \tag{4.3}
\end{array}
$$
\]

which are normalized such that

$$
\left\langle\psi_{i} \mid \psi_{j}\right\rangle=\psi_{i}^{\dagger} \psi_{j}=\delta_{i j}, \quad i, j=1,2,3,4 .
$$

The corresponding normalization factor $N$ obtained is

$$
N=\sqrt{\frac{W+\mu}{2 W}}, \quad \text { where } \quad W=+\left(p^{2}+\mu^{2}\right)^{1 / 2}, \mu=m c
$$

The solutions $\psi_{1}$ and $\psi_{2}$ correspond to positive energy states $E=+\left(p^{2} c^{2}+\right.$ $\left.m^{2} c^{4}\right)^{1 / 2}=+c W$ and the solutions $\psi_{3}$ and $\psi_{4}$ correspond to negative energy states with $E=-c W$.

If we assume the neutrino to be propagated along the $z$ direction and its rest mass to be zero, then $\mu=m c=0, W=p_{z}, p_{+}=p_{-}=0$. Substituting these values in Eqs. (4.2) and (4.3), we obtain the four solutions of the Dirac equation for the neutrino.

$$
\begin{align*}
& \psi_{1}=\frac{1}{\sqrt{2}}\left[\begin{array}{l}
1 \\
0 \\
1 \\
0
\end{array}\right]=\frac{1}{\sqrt{2}}\left[\begin{array}{c}
\chi_{+} \\
\chi_{+}
\end{array}\right], \psi_{2}=\frac{1}{\sqrt{2}}\left[\begin{array}{r}
0 \\
1 \\
0 \\
-1
\end{array}\right]=\frac{1}{\sqrt{2}}\left[\begin{array}{c}
\chi_{-} \\
-\chi_{-}
\end{array}\right],  \tag{4.4}\\
& \psi_{3}=\frac{1}{\sqrt{2}}\left[\begin{array}{r}
-1 \\
0 \\
1 \\
0
\end{array}\right]=\frac{1}{\sqrt{2}}\left[\begin{array}{r}
-\chi_{+} \\
\chi_{+}
\end{array}\right], \quad \psi_{4}=\frac{1}{\sqrt{2}}\left[\begin{array}{l}
0 \\
1 \\
0 \\
1
\end{array}\right]=\frac{1}{\sqrt{2}}\left[\begin{array}{c}
\chi_{-} \\
\chi_{-}
\end{array}\right] . \tag{4.5}
\end{align*}
$$

In Eqs. (4.4) and (4.5), $\chi_{+}$and $\chi_{-}$denote the two-component spinors

$$
\chi_{+}=\left[\begin{array}{l}
1 \\
0
\end{array}\right] ; \quad \chi_{-}=\left[\begin{array}{l}
0 \\
1
\end{array}\right] .
$$

The solutions (4.4) and (4.5) are obtained using the standard (Dirac) representation for the Dirac matrices $\boldsymbol{\alpha}$ and $\beta$.

$$
\boldsymbol{\alpha}=\left[\begin{array}{ll}
0 & \boldsymbol{\sigma}  \tag{4.6}\\
\boldsymbol{\sigma} & 0
\end{array}\right], \quad \beta=\left[\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right] .
$$

In addition, we can define the following gamma matrices.

$$
\begin{array}{ll}
\gamma_{0}=\beta, & \gamma=\beta \boldsymbol{\alpha}=\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \\
-\boldsymbol{\sigma} & 0
\end{array}\right], \\
\gamma_{5}=\gamma_{0} \gamma_{x} \gamma_{y} \gamma_{z}=-i\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right], & \gamma_{5}^{\prime}=i \gamma_{5}=\left[\begin{array}{cc}
0 & 1 \\
1 & 0
\end{array}\right] . \tag{4.7}
\end{array}
$$

The states $\psi_{1}$ and $\psi_{2}$ correspond to positive energy and the states $\psi_{3}$ and $\psi_{4}$ correspond to negative energy. It is found that the states $\psi_{1}, \psi_{2}, \psi_{3}, \psi_{4}$ are also the eigenfunctions of the chirality operator $\gamma_{5}^{\prime}$ and the helicity operator $\boldsymbol{\sigma} \cdot \hat{\boldsymbol{p}}$.

$$
\begin{array}{cl}
\gamma_{5}^{\prime} \psi_{1}=+\psi_{1}, & \text { eigenvalue of } \gamma_{5}^{\prime}:+1, \\
\gamma_{5}^{\prime} \psi_{2}=-\psi_{2}, & \text { eigenvalue of } \gamma_{5}^{\prime}:-1, \\
\gamma_{5}^{\prime} \psi_{3}=-\psi_{3}, & \text { eigenvalue of } \gamma_{5}^{\prime}:-1, \\
\gamma_{5}^{\prime} \psi_{4}=+\psi_{4}, & \text { eigenvalue of } \gamma_{5}^{\prime}:+1 . \\
\boldsymbol{\sigma} \cdot \hat{\boldsymbol{p}} \psi_{1}=+\psi_{1}, & \text { helicity }:+1, \\
\boldsymbol{\sigma} \cdot \hat{\boldsymbol{p}} \psi_{2}=-\psi_{2}, & \text { helicity }:-1, \\
\boldsymbol{\sigma} \cdot \hat{\boldsymbol{p}} \psi_{3}=+\psi_{3}, & \text { helicity }:+1,  \tag{4.9}\\
\boldsymbol{\sigma} \cdot \hat{\boldsymbol{p}} \psi_{4}=-\psi_{4}, & \text { helicity }:-1 .
\end{array}
$$

Thus, we find two positive energy states, one with positive helicity and the other with negative helicity. Similarly we have two negative energy states, one with positive helicity and the other with negative helicity. Experimentally, it is found that we have only neutrinos with negative helicity and anti-neutrinos with positive helicity. In other words, the neutrino is found to have its spin $\boldsymbol{S}$ oriented opposite to the direction of its momentum $\boldsymbol{p}$ whereas the anti-neutrino has its spin $\boldsymbol{S}$ oriented along the direction of its momentum $\boldsymbol{p}$, as shown in Fig. 4.1. This can


Figure 4.1 be realized by means of the projection operator $\frac{1}{2}\left(1-\gamma_{5}^{\prime}\right)$.

$$
\begin{array}{ll}
\frac{1}{2}\left(1-\gamma_{5}^{\prime}\right) \psi_{1}=0, & \frac{1}{2}\left(1-\gamma_{5}^{\prime}\right) \psi_{2}=\psi_{2}, \\
\frac{1}{2}\left(1-\gamma_{5}^{\prime}\right) \psi_{3}=\psi_{3}, & \frac{1}{2}\left(1-\gamma_{5}^{\prime}\right) \psi_{4}=0 \tag{4.10}
\end{array}
$$

The projection operator picks out the left-handed neutrino (neutrino with negative helicity) and right-handed anti-neutrino (anti-neutrino with positive helicity).

### 4.1.1 $\quad$ V-A form for the Weak interaction

In Fermi's theory of beta decay,

$$
n \rightarrow p+e^{-}+\bar{\nu}_{e}
$$

the decay rate is given by the product of nucleonic and leptonic current.

$$
j_{\mu}^{N} j_{\mu}^{\ell}=\left(\tilde{\psi}_{p} \gamma_{\mu} \psi_{n}\right)\left(\tilde{\psi}_{e} \gamma_{\mu} \psi_{\bar{\nu}}\right)
$$

The lepton current is of the form

$$
\begin{align*}
j_{\mu}^{\ell} & =\tilde{\psi}_{e} \gamma_{\mu} \psi_{\bar{\nu}} \\
& =\frac{1}{2} \tilde{\psi}_{e} \gamma_{\mu}\left(1-\gamma_{5}^{\prime}\right) \psi_{3} \tag{4.11}
\end{align*}
$$

The projection operator which eliminated the left-handed anti-neutrino gives the vector-axial vector current form (vector $\gamma_{\mu}-$ Axial vector $\gamma_{\mu} \gamma_{5}^{\prime}$ form) for the beta decay interaction. Such a form was originally proposed by Sudarshan and Marshak ${ }^{3}$ and by Feynman and Gell-Mann ${ }^{4}$. Here, we find such a form emerges from the experimental observation that the neutrinos are left-handed and the antineutrinos are right-handed.

In a similar way, we can modify the nucleonic current

$$
j_{\mu}^{N}=\tilde{\psi}_{p} \gamma_{\mu} \psi_{n}
$$

which is a pure vector current as given by Fermi, to include also the axial vector part.

$$
\begin{equation*}
j_{\mu}^{N}=\tilde{\psi}_{p} \gamma_{\mu}\left(1-\gamma_{5}^{\prime}\right) \psi_{n} \tag{4.12}
\end{equation*}
$$

The matrix element for the beta decay is obtained by coupling of nucleonic and leptonic current.

$$
\begin{equation*}
M=\frac{G_{F}}{\sqrt{2}}\left\{\tilde{\psi}_{p} \gamma_{\mu}\left(1-\gamma_{5}^{\prime}\right) \psi_{n}\right\}\left\{\tilde{\psi}_{e} \gamma_{\mu}\left(1-\gamma_{5}^{\prime}\right) \psi_{\bar{\nu}}\right\} \tag{4.13}
\end{equation*}
$$

where $G_{F}$ is the Fermi coupling constant.

[^24]
### 4.1.2 Weyl's two-component theory of neutrino

If you assume the neutrino mass to be zero, then the Dirac equation (4.1) for the neutrino becomes

$$
\begin{equation*}
c \boldsymbol{\alpha} \cdot \boldsymbol{p} \psi=E \psi \tag{4.14}
\end{equation*}
$$

The Dirac equation for a fermion with mass involves four anticommuting Dirac matrices $\alpha_{x}, \alpha_{y}, \alpha_{z}, \beta$ and to find a suitable representation for them, we require a minimum of $4 \times 4$ dimensional matrices. In the case of neutrino, since the mass term is absent, we have only three anticommuting matrices and so they can be represented by the three Pauli matrices $\sigma_{x}, \sigma_{y}, \sigma_{z}$ of dimension $2 \times 2$. So, the Dirac equation for the neutrino can be written as

$$
\begin{equation*}
c \boldsymbol{\sigma} \cdot \boldsymbol{p} \phi=E \phi \tag{4.15}
\end{equation*}
$$

where $\phi$ denotes a two-component wave function (spinor). Operating on the left by $\boldsymbol{\sigma} \cdot \boldsymbol{p}$ on both sides of the above equation and using the relation $(\boldsymbol{\sigma} \cdot \boldsymbol{p})(\boldsymbol{\sigma} \cdot \boldsymbol{p})=p^{2}$, we get

$$
c p^{2} \phi=E \boldsymbol{\sigma} \cdot \boldsymbol{p} \phi=\left(E^{2} / c\right) \phi .
$$

Equivalently

$$
E^{2}=p^{2} c^{2} \quad \text { or } \quad E= \pm p c .
$$

Substituting $E=+p c$ and $E=-p c$, in turn in Eq. (4.15), we get

$$
\begin{align*}
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{p} \phi_{+} & =\phi_{+}  \tag{4.16}\\
-\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{p} \phi_{-} & =\phi_{-} . \tag{4.17}
\end{align*}
$$

If we take the neutrino momentum to be along the z-direction, then $\phi_{+}$ and $\phi_{-}$correspond to the two helicity states $\chi_{+}$and $\chi_{-}$.

$$
\chi_{+}=\left[\begin{array}{l}
1 \\
0
\end{array}\right] ; \quad \chi_{-}=\left[\begin{array}{l}
0 \\
1
\end{array}\right] .
$$

Eqs. (4.16) and (4.17) can be physically interpreted to represent the positive helicity state for the neutrino and negative helicity state for antineutrino. But experimentally, it is found the other way. This should cause no problem. We have chosen the three Pauli matrices $\sigma_{x}, \sigma_{y}, \sigma_{z}$ as the three anticommuting matrices in Eq. (4.15). Instead, we can choose
$-\sigma_{x},-\sigma_{y},-\sigma_{z}$ as the three anticommuting matrices and we can rewrite Eq. (4.15) as

$$
\begin{equation*}
-c \boldsymbol{\sigma} \cdot \boldsymbol{p} \phi=E \phi, \tag{4.18}
\end{equation*}
$$

Since the sign is now reversed, we will obtain the correct helicity for the neutrino and anti-neutrino as observed experimentally.

$$
\begin{align*}
-\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{p} \phi_{-} & =\phi_{-}  \tag{4.19}\\
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{p} \phi_{+} & =\phi_{+} . \tag{4.20}
\end{align*}
$$

Eq. (4.19) represents the left-handed neutrino and Eq. (4.20) represents the right-handed antineutrino.

Although the two-component theory for neutrino is quite sufficient, we often study the interaction of neutrinos with other Fermions with mass, which can be represented only by four-component wave functions. In order to have a unified description in such cases, we need to introduce an appropriate four-component (bispinor) wave function for the neutrino. This can be done by using the Weyl representation for the Dirac matrices $\alpha$ and $\beta$.

### 4.1.3 The Weyl representation for Dirac matrices

Let us consider the Dirac equation given by Eq. (4.1). The Dirac matrices $\alpha_{x}, \alpha_{y}, \alpha_{z}$ and $\beta$ are anticommuting and each of which, when squared, yields unity.

$$
\alpha_{i} \alpha_{j}+\alpha_{j} \alpha_{i}=2 \delta_{i j}, \quad \alpha_{i} \beta+\beta \alpha_{i}=2 \delta_{\alpha_{i}, \beta}, \quad i, j=x, y, z
$$

It is possible to have more than one representation for the matrices, obeying the above properties. In Eq. (4.7), we have one representation, known as the Dirac representation which is also known as the standard representation. Here we shall use another representation, known as the Weyl or the chiral representation.

$$
\boldsymbol{\alpha}_{W}=\left[\begin{array}{rr}
\boldsymbol{\sigma} & 0  \tag{4.21}\\
0 & -\boldsymbol{\sigma}
\end{array}\right], \quad \beta_{W}=\left(\gamma_{0}\right)_{W}=\left[\begin{array}{rr}
0 & -1 \\
-1 & 0
\end{array}\right] .
$$

One can go from the Dirac representation to the Weyl representation by means of unitary transformation $U$.

$$
U=\frac{1}{\sqrt{2}}\left[\begin{array}{rr}
1 & 1  \tag{4.22}\\
-1 & 1
\end{array}\right]
$$

It can be easily verified that

$$
\begin{equation*}
\alpha_{W}=U \alpha U^{\dagger} ; \quad \beta_{W}=U \beta U^{\dagger} ; \tag{4.23}
\end{equation*}
$$

where $\alpha$ and $\beta$ without any suffix denote the Dirac matrices in the Dirac or standard representation. It can also be checked that the matrices $\gamma\left(\gamma_{x}, \gamma_{y}, \gamma_{z}\right)$ are the same in both the Dirac and Weyl representation.

$$
\gamma=\beta \boldsymbol{\alpha}=\beta_{W} \boldsymbol{\alpha}_{W}=\left[\begin{array}{rr}
0 & \boldsymbol{\sigma}  \tag{4.24}\\
-\boldsymbol{\sigma} & 0
\end{array}\right] .
$$

But the product of matrices $\gamma_{5}^{\prime}=i \gamma_{0} \gamma_{x} \gamma_{y} \gamma_{z}$ is different in the two representations.

$$
\begin{array}{ll}
\text { Dirac representation } & : \quad \gamma_{5}^{\prime}=\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right] \\
\text { Weyl representation } & :\left(\gamma_{5}^{\prime}\right)_{W}=\left[\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right] \tag{4.26}
\end{array}
$$

In the Weyl representation, the Dirac equation (4.1) becomes

$$
\left\{c\left[\begin{array}{cc}
\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0  \tag{4.27}\\
0 & -\boldsymbol{\sigma} \cdot \boldsymbol{p}
\end{array}\right]+\left[\begin{array}{cc}
0 & -m c^{2} \\
-m c^{2} & 0
\end{array}\right]\right\}\left[\begin{array}{l}
\phi_{+} \\
\phi_{-}
\end{array}\right]=E\left[\begin{array}{l}
\phi_{+} \\
\phi_{-}
\end{array}\right],
$$

where the Dirac bi-spinor $\psi$ is written in the two-component form, $\psi=$ $\left[\begin{array}{c}\phi_{+} \\ \phi_{-}\end{array}\right]$. Equation (4.27) leads to two coupled equations in $\phi_{+}$and $\phi_{-}$.

$$
\begin{align*}
c \boldsymbol{\sigma} \cdot \boldsymbol{p} \phi_{+}-m c^{2} \phi_{-} & =E \phi_{+}  \tag{4.28}\\
-c \boldsymbol{\sigma} \cdot \boldsymbol{p} \phi_{-}-m c^{2} \phi_{+} & =E \phi_{-} \tag{4.29}
\end{align*}
$$

If $m=0$, the Dirac equation splits into two separate equations in twocomponent spinors $\phi_{+}$and $\phi_{-}$.

$$
\begin{aligned}
c \boldsymbol{\sigma} \cdot \boldsymbol{p} \phi_{+} & =E \phi_{+}, \quad E= \pm c p . \\
-c \boldsymbol{\sigma} \cdot \boldsymbol{p} \phi_{-} & =E \phi_{-}, \quad
\end{aligned}
$$

$\underline{E=+c p}$

$$
\begin{align*}
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{p} \phi_{+} & =\phi_{+} ; & & \text {(Right-handed neutrino) }  \tag{4.31}\\
-\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{p} \phi_{-} & =\phi_{-} . & & \text {(Left-handed neutrino) } \tag{4.32}
\end{align*}
$$

$\underline{E=-c p}$

$$
\begin{align*}
-\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{p} \phi_{+} & =\phi_{+} ; & & \text {(Left-handed antineutrino) }  \tag{4.33}\\
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{p} \phi_{-} & =\phi_{-} ; & & \text {(Right-handed antineutrino) } \tag{4.34}
\end{align*}
$$

Since it is found experimentally that the neutrino is left-handed and the antineutrino is right-handed, Eqs. (4.32) and (4.34) alone are valid equations. So, in the four-component formalism, we should device a method by which the other two components vanish by using projection operators. In the Weyl representation ${ }^{5}$, we have

$$
\begin{align*}
\gamma & =\beta \alpha=\left[\begin{array}{rr}
0 & -1 \\
-1 & 0
\end{array}\right]\left[\begin{array}{rr}
\boldsymbol{\sigma} & 0 \\
0 & -\boldsymbol{\sigma}
\end{array}\right]=\left[\begin{array}{cc}
0 & \boldsymbol{\sigma} \\
-\boldsymbol{\sigma} & 0
\end{array}\right]  \tag{4.35}\\
\gamma_{5}^{\prime} & =i \gamma_{0} \gamma_{x} \gamma_{y} \gamma_{z}=\left[\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right] . \tag{4.36}
\end{align*}
$$

If $\psi$ is the four-component Dirac wave function for the mass zero particle, then the neutrino and antineutrino wave functions $\psi_{-}$and $\psi_{+}$can be obtained using the projection operators $\frac{1}{2}\left(1-\gamma_{5}^{\prime}\right)$ and $\frac{1}{2}\left(1+\gamma_{5}^{\prime}\right)$.

$$
\begin{align*}
& \psi_{-}=\frac{1}{2}\left(1-\gamma_{5}^{\prime}\right) \psi=\left[\begin{array}{c}
0 \\
\phi_{-}
\end{array}\right] .  \tag{4.37}\\
& \psi_{+}=\frac{1}{2}\left(1+\gamma_{5}^{\prime}\right) \psi=\left[\begin{array}{c}
\phi_{+} \\
0
\end{array}\right] . \tag{4.38}
\end{align*}
$$

Also, it can be verified that $\gamma_{5}^{\prime 2}=1$ and

$$
\begin{align*}
\gamma_{5}^{\prime} \psi_{-} & =-\psi_{-} .  \tag{4.39}\\
\gamma_{5}^{\prime} \psi_{+} & =\psi_{+} . \tag{4.40}
\end{align*}
$$

Hence, $\psi_{-}$and $\psi_{+}$are the solutions of the Dirac equation for the particle with zero rest mass and they are the eigenfunctions of the Hamiltonian, the helicity operator and of $\gamma_{5}^{\prime}$. Thus, we have shown that the two-component Weyl theory is equivalent to the four-component Dirac theory for the neutrino, if we choose the Weyl representation for the Dirac matrices $\boldsymbol{\alpha}$ and $\beta$.

[^25]
### 4.2 The Majorana neutrino

The Dirac theory of the electron suggested that for every particle, there is an antiparticle. In some cases, the particle and its antiparticle are identical as in the case of photon and $\pi^{0}$ meson. Majorana suggested that such an identity may exist in the case of neutrino too.

One goes from a particle to its antiparticle state by means of a charge conjugation operator $C=i \gamma_{y} \gamma_{0}$.

$$
\begin{align*}
\psi_{c}(x, t) & =C \gamma_{0} \psi^{*}(x, t)=i \gamma_{y} \gamma_{0} \gamma_{0} \psi^{*}(x, t) \\
& =i \gamma_{y} \psi^{*}(x, t), \quad \text { since } \gamma_{0}^{2}=1 \\
& =i\left[\begin{array}{cc}
0 & \sigma_{y} \\
-\sigma_{y} & 0
\end{array}\right] \psi^{*}(x, t) . \tag{4.41}
\end{align*}
$$

Let us write down explicitly the particle state and obtain its antiparticle state by the application of the charge conjugation operator.

$$
\psi=\left[\begin{array}{l}
\phi_{1}  \tag{4.42}\\
\phi_{2} \\
\phi_{3} \\
\phi_{4}
\end{array}\right]=\left[\begin{array}{l}
\phi_{+} \\
\phi_{-}
\end{array}\right],
$$

where

$$
\phi_{+}=\left[\begin{array}{c}
\phi_{1} \\
\phi_{2}
\end{array}\right] ; \quad \phi_{-}=\left[\begin{array}{c}
\phi_{3} \\
\phi_{4}
\end{array}\right] .
$$

The charge conjugated state $\psi_{c}$ is

$$
\begin{align*}
\psi_{c} & =i\left[\begin{array}{cc}
0 & \sigma_{y} \\
-\sigma_{y} & 0
\end{array}\right]\left[\begin{array}{c}
\phi_{+} \\
\phi_{-}
\end{array}\right]^{*} \\
& =i\left[\begin{array}{r}
\sigma_{y} \phi_{-}^{*} \\
-\sigma_{y} \phi_{+}^{*}
\end{array}\right] . \tag{4.43}
\end{align*}
$$

Majorana's proposal that the neutrino and the antineutrino are one and the same amounts to $\psi_{c}=\psi$.

$$
i\left[\begin{array}{r}
\sigma_{y} \phi_{-}^{*}  \tag{4.44}\\
-\sigma_{y} \phi_{+}^{*}
\end{array}\right]=\left[\begin{array}{c}
\phi_{+} \\
\phi_{-}
\end{array}\right] .
$$

Writing Eq. (4.44) in the four component form, we obtain the condition for the neutrino to be self-conjugate.

$$
\left[\begin{array}{r}
\phi_{4}^{*}  \tag{4.45}\\
-\phi_{3}^{*} \\
-\phi_{2}^{*} \\
\phi_{1}^{*}
\end{array}\right]=\left[\begin{array}{l}
\phi_{1} \\
\phi_{2} \\
\phi_{3} \\
\phi_{4}
\end{array}\right] .
$$

Thus we find the condition for $\psi$ to be self-conjugate is $\phi_{1}=\phi_{4}^{*}$ and $\phi_{2}=-\phi_{3}^{*}$. From an inspection of Eq. (4.44), we get the equivalent condition $\phi_{+}=i \sigma_{y} \phi_{-}^{*}$ for $\psi$ to be self-conjugate. Thus

$$
\psi=\left[\begin{array}{c}
i \sigma_{y} \phi_{-}^{*}  \tag{4.46}\\
\phi_{-}
\end{array}\right]=\left[\begin{array}{l}
\phi_{c} \\
\phi_{-}
\end{array}\right],
$$

where we have used a simplified notation $\phi_{c}$ for $i \sigma_{y} \phi_{-}^{*}$.
Now let us consider the four-component Dirac equation for the mass zero particle in the Wyle representation.

$$
\begin{align*}
c \boldsymbol{\alpha} \cdot \boldsymbol{p} \psi & =E \psi  \tag{4.47}\\
c\left[\begin{array}{cc}
\boldsymbol{\sigma} \cdot \boldsymbol{p} & 0 \\
0 & -\boldsymbol{\sigma} \cdot \boldsymbol{p}
\end{array}\right]\left[\begin{array}{c}
\phi_{+} \\
\phi_{-}
\end{array}\right] & =E\left[\begin{array}{l}
\phi_{+} \\
\phi_{-}
\end{array}\right] .
\end{align*}
$$

This yields two independent equations.

$$
\begin{align*}
c \boldsymbol{\sigma} \cdot \boldsymbol{p} \phi_{+} & =E \phi_{+} .  \tag{4.48}\\
-c \boldsymbol{\sigma} \cdot \boldsymbol{p} \phi_{-} & =E \phi_{-} . \tag{4.49}
\end{align*}
$$

If the neutrino is a Majorana particle, then $\phi_{+}=\phi_{c}=i \sigma_{y} \phi_{-}^{*}$ and the above two equations become

$$
\begin{align*}
c \boldsymbol{\sigma} \cdot \boldsymbol{p} \phi_{c} & =E \phi_{c} ; \text { where } \phi_{c}=i \sigma_{2} \phi_{-}^{*} .  \tag{4.50}\\
-c \boldsymbol{\sigma} \cdot \boldsymbol{p} \phi_{-} & =E \phi_{-} . \tag{4.51}
\end{align*}
$$

In the present case, the two equations are not independent and one can obtain Eq. (4.50) from (4.51). Let us start with Eq. (4.51) and take its complex conjugate.

$$
\begin{align*}
-c\left(\sigma_{x} p_{x}+\sigma_{y} p_{y}+\sigma_{z} p_{z}\right)^{*} \phi_{-}^{*} & =E \phi_{-}^{*} \\
-c\left(\sigma_{x} p_{x}-\sigma_{y} p_{y}+\sigma_{z} p_{z}\right) \phi_{-}^{*} & =E \phi_{-}^{*} \tag{4.52}
\end{align*}
$$

In Eq. (4.52), $\sigma_{x}, \sigma_{z}$ are real matrices and only $\sigma_{y}$ is purely imaginary. So, $\sigma_{y}^{*}=-\sigma_{y}$. Multiplying on the left by $i \sigma_{y}$ on both sides of Eq. (4.52), and remembering that $\sigma_{y}$ anticommutes with $\sigma_{x}$ and $\sigma_{z}$, we get

$$
\begin{align*}
-i c \sigma_{y}\left(\sigma_{x} p_{x}-\sigma_{y} p_{y}+\sigma_{z} p_{z}\right) \phi_{-}^{*} & =E\left(i \sigma_{y} \phi_{-}^{*}\right) \\
c\left(\sigma_{x} p_{x}+\sigma_{y} p_{y}+\sigma_{z} p_{z}\right)\left(i \sigma_{y} \phi_{-}^{*}\right) & =E\left(i \sigma_{y} \phi_{-}^{*}\right) \\
c(\boldsymbol{\sigma} \cdot \boldsymbol{p}) \phi_{c} & =E \phi_{c} . \tag{4.53}
\end{align*}
$$

Eq. (4.53) is identical with Eq. (4.50) and we have shown that Eqs. (4.51) and (4.50) are one and the same.

In the two-component theory, we had both left and right-handed neutrinos. Here we have only left-handed neutrino field $\phi_{-}$and the righthanded neutrino field $\phi_{c}$ is obtained from $\phi_{-}$. What is the physical significance of considering the neutrino as a Majorana particle?

If we consider the Dirac Lagrangian density $\mathscr{L}=\bar{\psi}\left(i \gamma^{\mu} \partial_{\mu}-m\right) \psi$, the mass term involves $\bar{\psi} \psi$. If $\bar{\psi} \psi$ vanishes, a Dirac particle cannot have any mass.

It is always possible to express $\psi$ as a sum of positive and negative helicity states $\psi_{R}$ and $\psi_{L}$, where $\psi_{R, L}=\frac{1}{2}\left(1 \pm \gamma_{5}^{\prime}\right) \psi$.

$$
\begin{align*}
\bar{\psi} \psi & =\left(\bar{\psi}_{R}+\bar{\psi}_{L}\right)\left(\psi_{R}+\psi_{L}\right) \\
& =\bar{\psi}_{R} \psi_{R}+\bar{\psi}_{L} \psi_{L}+\bar{\psi}_{R} \psi_{L}+\bar{\psi}_{L} \psi_{R} \\
& =\bar{\psi}_{R} \psi_{L}+\bar{\psi}_{L} \psi_{R}, \tag{4.54}
\end{align*}
$$

since $\bar{\psi}_{R} \psi_{R}=\bar{\psi}_{L} \psi_{L}=0$, as shown below.

$$
\begin{align*}
\bar{\psi}_{R} \psi_{R} & =\frac{1}{4}\left\{\left(1+\gamma_{5}^{\prime}\right) \psi\right\}^{\dagger} \gamma_{0}\left\{\left(1+\gamma_{5}^{\prime}\right) \psi\right\} \\
& =\frac{1}{4} \psi^{\dagger}\left(1+\gamma_{5}^{\prime}\right) \gamma_{0}\left(1+\gamma_{5}^{\prime}\right) \psi \\
& =\frac{1}{4} \bar{\psi}\left(1-\gamma_{5}^{\prime}\right)\left(1+\gamma_{5}^{\prime}\right) \psi \\
& =\bar{\psi}_{L} \psi_{L}=0 . \tag{4.55}
\end{align*}
$$

In deducing the relation (4.55), we have used the following properties of $\gamma_{5}^{\prime}$ :

$$
\gamma_{5}^{\prime \dagger}=\gamma_{5}^{\prime} ; \quad \gamma_{5}^{\prime} \gamma_{0}=-\gamma_{0} \gamma_{5}^{\prime} ; \quad \gamma_{5}^{\prime 2}
$$

Thus, we find that a mass term for the Dirac field can be introduced only if there are $\psi_{L}$ and $\psi_{R}$ chiral components. Since the right-handed neutrinos do not exist, it is not possible for neutrinos to have Dirac mass. In the case of Majorana neutrinos, with the self-conjugate field, the mass term is given by

$$
\begin{align*}
\bar{\psi} \psi=\psi^{\dagger} \gamma_{0} \psi & =\left[\begin{array}{ll}
-i \phi_{-}^{T} \sigma_{y} & \phi_{-}^{T^{*}}
\end{array}\right]\left[\begin{array}{rr}
0 & -1 \\
-1 & 0
\end{array}\right]\left[\begin{array}{c}
i \sigma_{y} \phi_{-}^{*} \\
\phi_{-}
\end{array}\right] \\
& =\left[\begin{array}{ll}
-i \phi_{-}^{T} \sigma_{y} & \phi_{-}^{T^{*}}
\end{array}\right]\left[\begin{array}{c}
-\phi_{-} \\
-i \sigma_{y} \phi_{-}^{*}
\end{array}\right] \\
& =\phi_{-}^{T}\left(i \sigma_{y}\right) \phi_{-}+\phi_{-}^{T^{*}}\left(-i \sigma_{y}\right) \phi_{-}^{*} \\
& =\phi_{-}^{T}\left(i \sigma_{y}\right) \phi_{-}+\text {h.c. } \tag{4.56}
\end{align*}
$$

since

$$
\left\{\phi_{-}^{T}\left(i \sigma_{y}\right) \phi_{-}\right\}^{\dagger}=\phi_{-}^{\dagger}\left(i \sigma_{y}\right)^{\dagger}\left(\phi_{-}^{T}\right)^{\dagger}=\phi_{-}^{T^{*}}\left(-i \sigma_{y}\right) \phi_{-}^{*} .
$$

The above analysis clearly shows that it is really possible to have mass for the single left-handed Majorana field $\phi_{-}$. Writing explicitly $\phi_{-}$as a two-component spinor $\left[\begin{array}{c}\phi_{3} \\ \phi_{4}\end{array}\right]$, we find

$$
\begin{equation*}
\phi_{-}^{T}\left(i \sigma_{y}\right) \phi_{-}=\phi_{3} \phi_{4}-\phi_{4} \phi_{3} . \tag{4.57}
\end{equation*}
$$

This implies that the components of the two-component spinor should be anticomputing. Further, unlike the Dirac mass term $\bar{\psi} \psi$ which is invariant under the global $\mathrm{U}(1)$ transformation, the Majorana mass term violates the global $U(1)$ symmetry, which is the lepton number in this case. The electron which is a spin- $\frac{1}{2}$ particle, carrying conserved $\mathrm{U}(1)$ charge cannot be a Majorana particle.

### 4.3 Neutrinoless double beta decay

Now, the question arises whether neutrino is a Dirac particle with zero rest mass or a Majorana particle with a possible rest mass. Is there a lepton number conservation or not? It is pointed out that the experimental observation of neutrinoless double beta decay ( $0 \nu \beta \beta$-decay)

$$
(A, Z) \rightarrow(A, Z+2)+e^{-}+e^{-}
$$

will settle this issue. Schechter and Valle ${ }^{6}$ showed that the observation of $0 \nu \beta \beta$-decay will ensure that the neutrinos are Majorana particles with non-zero rest mass, regardless of the mechanism that causes it. Also it will admit the violation of lepton number conservation.

We are all familiar with ordinary beta decay with the emission of single electron and a single neutrino.

$$
(A, Z) \rightarrow(A, Z+1)+e^{-}+\bar{\nu}_{e} .
$$

If the binding energy of $(\mathrm{A}, \mathrm{Z}+1)$ nucleus is less than that of $(\mathrm{A}, \mathrm{Z})$ nucleus but the binding energy of $(A, Z+2)$ is greater than that of $(A, Z)$, the single beta decay is not energetically possible but allows the double beta decay.

[^26]Under such favourable conditions, it is possible to observe two-neutrino double beta decay ( $2 \nu \beta \beta$ - decay).

$$
(A, Z) \rightarrow(A, Z+2)+e^{-}+e^{-}+\bar{\nu}_{e}+\bar{\nu}_{e}
$$

For the first time, in 1987, such a two-neutrino double-beta decay was observed by M. Moe and his collaborators ${ }^{7}$ with ${ }^{82}$ Se nucleus.

$$
{ }^{82} \mathrm{Se} \rightarrow{ }^{82} \mathrm{Kr}+e^{-}+e^{-}+\bar{\nu}_{e}+\bar{\nu}_{e} .
$$

Subsequently, many more $2 \nu \beta \beta$ were observed with other nuclei

$$
{ }^{48} \mathrm{Ca},{ }^{76} \mathrm{Ge},{ }^{96} \mathrm{Zr},{ }^{100} \mathrm{Mo},{ }^{116} \mathrm{Cd},{ }^{128} \mathrm{Te},{ }^{130} \mathrm{Te},{ }^{136} \mathrm{Xe},{ }^{150} \mathrm{Nd},{ }^{238} \mathrm{U} .
$$

The half-life of the parent nuclei in $2 \nu \beta \beta$-decay is of the order of $10^{21}$ years.

If the neutrino is a Majorana particle (for which neutrino and antineutrino are one and the same particle), it is possible that the neutrino emitted by one nucleon is absorbed by another nucleon in the nucleus and the neutrinoless double beta decay may occur.

There has been a claim ${ }^{8}$ of having observed the neutrinoless double beta decay of ${ }^{76} \mathrm{Ge}$ with a half-life $T_{1 / 2}=2.23_{-0.31}^{+0.44} \times 10^{25}$ but it has not yet been confirmed by others.

### 4.4 Three generations of leptons

The electron and the neutrino associated with it are called leptons. Subsequently, it was found that they are heavier particles known as muon ( $\mu$ with mass $105 \mathrm{MeV} / \mathrm{c}^{2}$ ) and tau ( $\tau$ with mass $1777 \mathrm{MeV} / \mathrm{c}^{2}$ ) particles and their associated neutrinos $\left(\nu_{\mu}, \nu_{\tau}\right)$ which exhibit similar interaction properties as electron and its neutrino $\left(\nu_{e}\right)$. Muons were first discovered in cosmic ray experiments around the year 1936 by C. D. Anderson, S. H. Neddermeyer and others and tau lepton was discovered in 1975 at the Stanford Linear Accelerator laboratory by M. L. Perl. There are known as three generations of leptons.


[^27]The absence of decay modes of muon into electron and gamma ray

$$
\mu^{-} \nrightarrow e^{-}+\gamma, \quad \mu^{+} \nrightarrow e^{+}+\gamma,
$$

although energetically possible and the absence of reactions

$$
\nu_{e}+n \nrightarrow \mu^{-}+p, \quad \nu_{\mu}+n \nrightarrow e^{-}+p,
$$

necessitated the introduction of separate lepton numbers for each generation of leptons. This has been fully supported by the study of other allowed reactions involving leptons.

$$
\begin{aligned}
n & \rightarrow p+e^{-}+\bar{\nu}_{e} . \\
\pi^{-} & \rightarrow \mu^{-}+\bar{\nu}_{\mu} . \\
\mu^{-} & \rightarrow e^{-}+\bar{\nu}_{e}+\nu_{\mu} . \\
\nu_{e}+n & \rightarrow e^{-}+p . \\
\nu_{\mu}+n & \rightarrow \mu^{-}+p .
\end{aligned}
$$

Thus, each generation of leptons is given a separate lepton number and it is found to be conserved in any interaction. This is incorporated in the Standard Model of elementary particles, according to which there are three different flavours of neutrino, each having zero rest mass. However, it is now known that the neutrino may change its flavour and the rest mass of the neutrino is not exactly zero but may be very small.

### 4.5 Neutrino oscillations

### 4.5.1 Experimental evidence

Neutrinos abound in nature. Thermo-nuclear reactions in the Sun produce low energy electron neutrinos in thousands of billions. Cosmic ray showers produce high energy electron and muon neutrinos in great measure. Manmade accelerators and nuclear reactors produce neutrinos of all flavours.

Since neutrino is a neutral particle without any charge, it cannot be observed directly but can only be detected indirectly by observing the reaction partners in the weak interaction, in which it participates. In nineteen fifties, Raymond Davis devised an ingenious experiment to detect the neutrinos using the reaction

$$
\begin{equation*}
\nu_{e}+{ }^{37} \mathrm{Cl} \rightarrow{ }^{37} \mathrm{Ar}+e^{-} . \tag{4.58}
\end{equation*}
$$

The Q -value for the reaction is 0.8 MeV and so the incident neutrino should have an energy greater than 0.8 MeV for the reaction to occur.

Since Davis has used the nuclear reactor as the source of neutrinos, he failed in his effort, since the nuclear reactor produced $\bar{\nu}_{e}$ rather than $\nu_{e}$. At the same time, Cowen and Reines used the reaction

$$
\bar{\nu}_{e}+p \rightarrow n+e^{+}
$$

for detecting the neutrino through its interaction with proton and they succeeded in their attempt by detecting the final products in the reaction without any ambiguity. Reines was rewarded with the Nobel prize for this achievement.

## Solar neutrinos

Undaunted by the failure to detect the neutrinos from the reactors, Davis and his collaborators continued their experiments with chlorine for the next two decades using solar neutrinos. The Sun serves as a natural source of electron neutrinos. Bahcal and his collaborators ${ }^{9}$ have developed a theory known as the Standard Solar Model (SSM) using the thermonuclear fusion reactions via the pp chain and the CNO cycle. This ultimately results in the fusion reaction

$$
4 p \rightarrow{ }^{4} \mathrm{He}+2 e^{+}+2 \nu_{e}+26.7 \mathrm{MeV},
$$

with a release of 26.7 MeV energy that serves as the source of solar energy. From this, one can estimate correctly the number of electron neutrinos received on the earth per second per sq. cm . There are about $7 \times 10^{9}$ neutrinos reaching the earth from the Sun per sq. cm. per sec. Most of the electron neutrinos emitted by the Sun are of low energy less than 0.5 MeV . So, only the high energy neutrinos, which form a small portion of the neutrino energy spectrum and come from the decay of Boron $\left({ }^{8} \mathrm{~B} \rightarrow\right.$ ${ }^{8} \mathrm{Be}+e^{+}+\nu_{e}$ ), will trigger the reaction (4.58).

## The Homestake Experiment

Davis used the Homestake gold mine in South Dakota (USA) for his solar neutrino experiments. A large tank containing 615 tons of fluid, rich in chlorine, called tetra-chloroethylene was used and the radioactive Argon atoms produced in the reaction (4.58) were extracted periodically from

[^28]the tank by a special technique. Since the neutrinos are weakly interacting particles, they produced argon atoms at a very slow rate of one argon atom per day. Every month the ${ }^{36} \mathrm{Ar}$ atoms were collected and counted. Instead of 30 atoms expected from the SSM calculations, only 10 were actually observed ${ }^{10}$. In other words, only one third of the expected number of neutrinos were counted in the experiment. This has come to be known as the Solar neutrino puzzle.

## The Kamiokande Experiment

After almost two decades, the deficiency in the solar neutrino count reported by the Homestack experiment, received an important confirmation from the Japanese Kamiokande experiment ${ }^{11}$, led by Masatoshi Koshiba and his team. The Kamiokande experiment observed only $40 \%$ of the expected solar neutrinos from the Sun.

The Kamiokande detector consisted of about 2 Kilotons of pure water, surrounded by an array of 1000 photomultiplier tubes. The solar neutrino undergoes scattering with an electron in the water molecule and the ejected high energy electron is detected by Cherenkov radiation it emits.

$$
\nu_{e}+e^{-} \rightarrow \nu_{e}+e^{-} .
$$

The energy threshold of Kamiokande detector is about 7 MeV and so only the high energy part of the Boron- 8 spectrum was observed. The deficiency observed in the solar neutrino count supported the solar-neutrino puzzle reported earlier by the Homestack experiment. There is a distinct difference between the two experiments. The Homestack experiment used a passive radioactive chemical method for detection of events whereas the Kamiokande experiment used an active method of registering the time at which the event took place and also the direction from which the neutrino came; thereby assuring that it came from the Sun. Thus the solar neutrino puzzle, first reported by the Homestack experiment was corroborated by the Kamiokande experiment. Both Raymond Davis and Masatoshi Koshiba were awarded the Nobel Prize in the year 2002 for having established a new field of research known as neutrino astronomy which is of great importance to elementary particle physics, astrophysics and cosmology.

[^29]
## The Sudbury Neutrino Observatory

Since the Standard Solar Model is found to be on a solid foundation and the deficiency in the count of electron neutrinos is also corroborated by other experimenters, the only option is to speculate a flavour change of electron neutrino to muon neutrino or tau neutrino during its flight from the Sun. The experiments done so far could not detect the neutrinos of other flavours. The Sudbury Neutrino observatory (SNO) located in a nickel mine in Ontario in Canada was designed to observe the neutrinos of all flavours in order to establish convincingly the phenomenon of neutrino oscillations.

In SNO, about 1000 tons of heavy water was used in a transparent sphere surrounded by 9,500 photomultiplier tubes. A neutrino interacting with the deuterium in the heavy water can induce the following two types of interactions - charged current and neutral current weak interactions.

$$
\begin{array}{ll}
\nu_{e}+d \rightarrow p+p+e^{-} & \text {charged current (CC) } \\
\nu_{x}+d \rightarrow p+n+\nu_{x}, \quad x=e, \mu, \tau & \text { neutral current (NC) } \tag{4.60}
\end{array}
$$

The energy threshold for these reactions is a few MeV and so only the solar neutrinos that come from ${ }^{8} \mathrm{~B}$ decay can only drive the aforesaid reactions. Only the electron neutrinos can induce the charged current reaction and the energetic electron that is produced in the CC reaction can be detected by the Cerenkov radiation that it emits. The muon neutrinos and tau neutrinos cannot induce CC reactions since the solar neutrinos do not have sufficient energy to produce their partners $\mu$ and $\tau$. On the other hand, the NC reaction can be induced by neutrinos of all flavours. The results ${ }^{12}$ of the CC reaction rate are in perfect agreement with the Homestack and Kamiokande experiments that only one third of the solar electron neutrinos are reaching the earth without any change of flavour. On the other hand, the NC reaction rate coincides with the rate predicted by Standard Solar Model of Bahcall. This can be considered as a clear evidence that two third of the electron neutrinos emitted by the Sun have changed their flavour during their passage to the earth.

## Atmospheric neutrinos

Just as the Sun acts as a rich source of electron neutrinos of low energy, cosmic rays offer as a rich source of high energy electron and muon neutrinos. Cosmic rays consist mostly of high energy protons which collide

[^30]with atmospheric atoms to produce pions and muons which in turn decay releasing high energy electron and muon neutrinos and antineutrinos. These are known as atmospheric neutrinos and they are in the energy range of GeV .
\[

$$
\begin{array}{ll}
\pi^{-} \rightarrow \mu^{-}+\bar{\nu}_{\mu} ; & \pi^{+} \rightarrow \mu^{+}+\nu_{\mu} \\
\mu^{-} \rightarrow e^{-}+\bar{\nu}_{e}+\nu_{\mu} ; & \mu^{+} \rightarrow e^{+}+\nu_{e}+\bar{\nu}_{\mu}
\end{array}
$$
\]

The ratio $R$ of the number of atmospheric muon neutrinos to the number of electron neutrinos is two.

$$
\begin{equation*}
R=\frac{\nu_{\mu}+\bar{\nu}_{\mu}}{\nu_{e}+\bar{\nu}_{e}}=2 . \tag{4.61}
\end{equation*}
$$

The experiments do not distinguish between the events caused by neutrinos and antineutrinos but differentiate between the events initiated by neutrinos of different flavours. So, in the discussion that follows, neutrinos include antineutrinos also. Besides the experiments have directional facility. They can count separately the number of neutrinos that come from above and also those that come from below.

## The Super-Kamiokande Experiment

The Japanese group working in Super-Kamiokande detector laboratory, located at a depth of 1000 meters below the ground level in a zinc mine, about 250 Km away from Tokyo in the north west direction, observed the electron and muon neutrinos that were produced by cosmic rays interacting with earth's atmosphere. It was found that the ratio $R=2$ for the neutrinos that came from above but this ratio $R$ got reduced to one for the neutrinos that came from below. A detailed analysis revealed that the number of electron neutrinos that came from below was almost the same as the number of electron neutrinos that came from above but the number of muon neutrinos that came from below was much less than the number of muon neutrinos that came from above. These results are quite perplexing. Why in the case of atmospheric electron neutrinos, there is no depletion but in the case of muon neutrinos, there is a depletion? Why, in the case of solar electron neutrinos, a depletion in the count was observed in solar experiments but there is no such depletion in the case of atmospheric electron neutrinos? You will find the answers in the solved problems 4.4, 4.5 and 4.6.

The muon neutrinos that came from below had to travel through the earth (earth's diameter $=12,742 \mathrm{Km}$ ) before reaching the detector. During this passage through the earth, some of the muon neutrinos might


Figure 4.2: Neutrinos from below should travel a distance of $12,742 \mathrm{Km}$ (diameter of the earth) before reaching the detector.
have changed into tau neutrinos and, as a consequence, the neutrino detector which could only detect muon neutrinos showed less count for the neutrinos that came from below when compared to the detector count for the neutrinos that came from above ${ }^{13}$. This was in the year 1998.

This is a clear evidence that the neutrinos, as they travel large distances change from one flavour to another and this has come to be known as neutrino oscillations. This is not possible if the neutrinos are massless. So, it is conjectured that the neutrinos carry a small mass. This neutrino oscillation was corroborated later in the year 2001 by the Sudbury Neutrino Observatory (SNO), located in Canada, using solar neutrinos.

The chief physicists Takaaki Kajita and Arther B. McDonald of these two large research groups - Super Kamiokande Collaboration and Sudbury Neutrino Observatory - were awarded the Nobel Prize in Physics for the year 2015.

### 4.5.2 Theoretical analysis

Neutrino is produced in weak interaction along with a charged lepton or is absorbed to produce a charged lepton. Since there are three types of charged leptons, electron, muon and tau, the neutrino, by definition, is a flovour eigenstate. The neutrino that is associated with electron is the electron neutrino, that associated with muon is the muon neutrino and that associated with tau is the tau neutrino.

Let us try to estimate the depletion of electron neutrinos during the

[^31]flight from the Sun by conversion to muon neutrinos and tau neutrinos. If the neutrinos have a mass, however small it may be, then we can define two sets of eigenstates, mass eigenstates $\left.\left|\nu_{1}\right\rangle,\left|\nu_{2}\right\rangle, \nu_{3}\right\rangle$ and weak interaction eigenstates $\left.\left|\nu_{e}\right\rangle,\left|\nu_{\mu}\right\rangle, \nu_{\tau}\right\rangle$. Let us assume that one can go from one set of eigenstates $\left|\nu_{i}\right\rangle$ to the other set $\left|\nu_{w}\right\rangle$ by unitary transformation.
\[

$$
\begin{equation*}
\left|\nu_{w}\right\rangle=\sum_{i=1}^{3} U_{w i}\left|\nu_{i}\right\rangle, \quad w=e, \mu, \tau \tag{4.62}
\end{equation*}
$$

\]

where $U_{w i}$ is the unitary transformation operator that obeys the following normalization conditions.

$$
\begin{equation*}
\sum_{w} U_{i w} U_{j w}^{*}=\delta_{i j} ; \quad \sum_{i} U_{w i} U_{w^{\prime} i}^{*}=\delta_{w w^{\prime}} . \tag{4.63}
\end{equation*}
$$

Let us assume that $H$ is diagonal in the $\left|\nu_{i}\right\rangle$ basis, such that

$$
\begin{equation*}
H(p)\left|\nu_{i}\right\rangle=E_{i}\left|\nu_{i}\right\rangle, \tag{4.64}
\end{equation*}
$$

where

$$
\begin{equation*}
E_{i}=\left(p^{2} c^{2}+m_{i}^{2} c^{4}\right)^{\frac{1}{2}} \approx p c\left\{1+\frac{m_{i}^{2} c^{2}}{p^{2}}\right\}^{\frac{1}{2}} \approx p c+\frac{1}{2} \frac{m_{i}^{2} c^{3}}{p} \tag{4.65}
\end{equation*}
$$

Let us now consider the time evolution of the mass eigenstate.

$$
\begin{equation*}
\left|\nu_{i}(t)\right\rangle=e^{-i E_{i} t / \hbar}\left|\nu_{i}(0)\right\rangle=e^{-i E_{i} t / \hbar}\left|\nu_{i}\right\rangle . \tag{4.66}
\end{equation*}
$$

The time evolution of the weak interaction eigenstate can be obtained from eq. (4.62).

$$
\begin{equation*}
\left|\nu_{w}(t)\right\rangle=\sum_{i=1}^{3} U_{w i} e^{-i E_{i} t / \hbar}\left|\nu_{i}\right\rangle, \quad w=e, \mu, \tau \tag{4.67}
\end{equation*}
$$

Now, we can write down the amplitude for the weak interaction amplitude $\left|\nu_{w}(t)\right\rangle$ at time $t$ to contain the weak interaction eigenstate $\left|\nu_{w^{\prime}}\right\rangle$.

$$
\begin{align*}
\left\langle\nu_{w^{\prime}} \mid \nu_{w}(t)\right\rangle & =\sum_{i=1}^{3} U_{w i} e^{-i E_{i} t / \hbar}\left\langle\nu_{w^{\prime}} \mid \nu_{i}\right\rangle \\
& =\sum_{i=1}^{3} U_{w i} e^{-i E_{i} t / \hbar} \sum_{j}\left\langle U_{w^{\prime} j} \nu_{j} \mid \nu_{i}\right\rangle \\
& =\sum_{i=1}^{3} U_{w i} e^{-i E_{i} t / \hbar} U_{w^{\prime} i}^{*} . \tag{4.68}
\end{align*}
$$

The last step is obtained using the orthonormality condition of the mass eigenstates $\left\langle\nu_{j} \mid \nu_{i}\right\rangle=\delta_{i j}$.

From Eq. (4.68), we obtain the probability $P_{w^{\prime} w}$ of finding the neutrino of type $w^{\prime}$ after time $t$.

$$
\begin{align*}
P_{w^{\prime} w} & =\left|\left\langle\nu_{w^{\prime}} \mid \nu_{w}(t)\right\rangle\right|^{2} \\
& =\left(\sum_{i=1}^{3} U_{w i} e^{-i E_{i} t / \hbar} U_{w^{\prime} i}^{*}\right)\left(\sum_{j=1}^{3} U_{w j} e^{-i E_{j} t / \hbar} U_{w^{\prime} j}^{*}\right)^{*} \\
& =\sum_{i, j}\left(U_{w i} U_{w^{\prime} i}^{*}\right)\left(U_{w j}^{*} U_{w^{\prime} j}\right)\left\{e^{-i\left(E_{i}-E_{j}\right) t / \hbar}-1+1\right\} \\
& =\delta_{w^{\prime} w}+\sum_{i, j} U_{w i} U_{w^{\prime} i}^{*} U_{w j}^{*} U_{w^{\prime} j}\left\{e^{-i\left(E_{i}-E_{j}\right) t / \hbar}-1\right\} . \tag{4.69}
\end{align*}
$$

Since the summation is over dummy indices $i$ and $j$, you are free to exchange them and include another term with $i$ and $j$ exchanged.

$$
\begin{align*}
P_{w^{\prime} w}= & \delta_{w^{\prime} w}+\frac{1}{2} \sum_{i, j}\left[U_{w i} U_{w^{\prime} i}^{*} U_{w j}^{*} U_{w^{\prime} j}\left\{e^{-i\left(E_{i}-E_{j}\right) t / \hbar}-1\right\}\right. \\
& \left.+U_{w j} U_{w^{\prime} j}^{*} U_{w i}^{*} U_{w^{\prime} i}\left\{e^{-i\left(E_{j}-E_{i}\right) t / \hbar}-1\right\}\right] . \tag{4.70}
\end{align*}
$$

If $U$ is complex, then it is possible to write ${ }^{14}$

$$
\begin{equation*}
U_{w i} U_{w^{\prime} i}^{*} U_{w j}^{*} U_{w^{\prime} j}=\left|U_{w i} U_{w^{\prime} i}^{*} U_{w j}^{*} U_{w^{\prime} j}\right| e^{i \phi} \tag{4.71}
\end{equation*}
$$

where $\phi$ is the overall phase factor. The exchange term is the complex conjugate of (4.71).

$$
\begin{equation*}
U_{w i}^{*} U_{w^{\prime} i} U_{w j} U_{w^{\prime} j}^{*}=\left|U_{w i} U_{w^{\prime} i}^{*} U_{w j}^{*} U_{w^{\prime} j}\right| e^{-i \phi} \tag{4.72}
\end{equation*}
$$

[^32]Substituting (4.71) and (4.72) into Eq. (4.70), we get

$$
\begin{align*}
P_{w^{\prime} w}= & \delta_{w^{\prime} w}+\frac{1}{2} \sum_{i, j}\left[\left|U_{w i} U_{w^{\prime} i}^{*} U_{w j}^{*} U_{w^{\prime} j}\right| e^{i \phi}\left\{e^{-i\left(E_{i}-E_{j}\right) t / \hbar}-1\right\}\right. \\
& \left.+\left|U_{w i} U_{w^{\prime} i}^{*} U_{w j}^{*} U_{w^{\prime} j}\right| e^{-i \phi}\left\{e^{i\left(E_{i}-E_{j}\right) t / \hbar}-1\right\}\right] \\
= & \delta_{w^{\prime} w}+\frac{1}{2} \sum_{i, j}\left|U_{w i} U_{w^{\prime} i}^{*} U_{w j}^{*} U_{w^{\prime} j}\right| \\
& \times\left\{e^{-i\left\{\left(E_{i}-E_{j}\right) t / \hbar-\phi\right\}}-e^{i \phi}+e^{i\left\{\left(E_{i}-E_{j}\right) t / \hbar-\phi\right\}}-e^{-i \phi}\right\} \\
= & \delta_{w^{\prime} w}+\frac{1}{2} \sum_{i, j}\left|U_{w i} U_{w^{\prime} i}^{*} U_{w j}^{*} U_{w^{\prime} j}\right| \\
& \times\left\{\left(2 \cos \frac{\left(E_{i}-E_{j}\right) t}{\hbar}-\phi\right)-2 \cos \phi\right\} \tag{4.73}
\end{align*}
$$

Let us, for simplicity, assume that $U$ is real. Then $\phi=0$ and Eq.(4.73) reduces to

$$
\begin{equation*}
P_{w^{\prime} w}=\delta_{w^{\prime} w}+\sum_{i, j} U_{w i} U_{w^{\prime} i} U_{w j} U_{w^{\prime} j}\left\{-2 \sin ^{2} \frac{\left(E_{i}-E_{j}\right) t}{2 \hbar}\right\} \tag{4.74}
\end{equation*}
$$

The quantity within the curly bracket in Eq. (4.74) can be expressed in terms of the neutrino masses using Eq. (4.65).

$$
\begin{align*}
E_{i}-E_{j} & =\left(p c+\frac{1}{2} \frac{m_{i}^{2} c^{3}}{p}\right)-\left(p c+\frac{1}{2} \frac{m_{j}^{2} c^{3}}{p}\right) \\
& =\frac{1}{2 p}\left(m_{i}^{2}-m_{j}^{2}\right) c^{3} . \tag{4.75}
\end{align*}
$$

If $L$ is the distance travelled by the neutrino, then the time taken for the neutrino to cover this distance $t \approx L / c$. Thus,

$$
\begin{equation*}
\frac{\left(E_{i}-E_{j}\right) t}{\hbar}=\frac{\left(E_{i}-E_{j}\right) L}{c \hbar}=\frac{\left(m_{i}^{2}-m_{j}^{2}\right) c^{3} L}{2 p c \hbar}=\frac{\left(m_{i}^{2}-m_{j}^{2}\right) c^{3} L}{2 E_{\nu} \hbar} \tag{4.76}
\end{equation*}
$$

where we have substituted $p c$ by $E_{\nu}$, since the neutrino mass is negligible.
The quantity within the curly bracket in Eq. (4.74) can now be expressed as

$$
\begin{equation*}
\{\cdots\}=-2 \sin ^{2} \frac{\left(E_{i}-E_{j}\right) t}{2 \hbar}=-2 \sin ^{2} \frac{\pi L}{\lambda_{i j}}, \tag{4.77}
\end{equation*}
$$

where

$$
\begin{equation*}
\lambda_{i j}=\frac{4 \pi(\hbar c) E_{\nu}}{\left(m_{i}^{2}-m_{j}^{2}\right) c^{4}} . \tag{4.78}
\end{equation*}
$$

Substituting (4.77) into Eq. (4.74), we obtain the probability of neutrino oscillation from flavour $w$ to flavour $w^{\prime}$ for a distance of travel $L$.

$$
\begin{align*}
P_{w^{\prime} w} & =\delta_{w^{\prime} w}-2 \sum_{i, j} U_{w i} U_{w^{\prime} i} U_{w j} U_{w^{\prime} j} \sin ^{2}\left(\frac{\pi L}{\lambda_{i j}}\right) \\
& =\delta_{w^{\prime} w}-4 \sum_{i<j} U_{w i} U_{w^{\prime} i} U_{w j} U_{w^{\prime} j} \sin ^{2}\left(\frac{\pi L}{\lambda_{i j}}\right) . \tag{4.79}
\end{align*}
$$

The oscillation is permissible only if the mass eigenstates are non-degenerate $\left(m_{i} \neq m_{j}\right)$. By substituting the values of the physical constants:

$$
\hbar=1.055 \times 10^{-34} \mathrm{Js} ; \quad c=2.998 \times 10^{8} \mathrm{~m} / \mathrm{s} ; \quad 1 \mathrm{eV}=1.602 \times 10^{-19} \mathrm{~J}
$$

the oscillation length parameter $\lambda_{i j}$ and its dependence on the neutrino energy and their mass differences can be explicitly given.

$$
\begin{align*}
\hbar c & =1.055 \times 2.998 \times 10^{-26} \mathrm{Jm} \\
& =\frac{1.055 \times 2.998 \times 10^{-26}}{1.602 \times 10^{-19}}=1.9743 \times 10^{-7} \mathrm{eV} \mathrm{~m} \\
\lambda_{i j} & =\frac{4 \pi\left(1.9743 \times 10^{-7}\right) E_{\nu}}{\left(m_{i}^{2}-m_{j}^{2}\right) c^{4}} \\
& =\frac{4 \pi\left(1.9743 \times 10^{-1}\right)}{\left(m_{i}^{2}-m_{j}^{2}\right) c^{4}}\left(\frac{E_{\nu}}{10^{6}}\right) \\
& =2.481\left(\frac{E_{\nu}}{\mathrm{MeV}}\right) \frac{(\mathrm{eV})^{2}}{\left(m_{i}^{2}-m_{j}^{2}\right)} \tag{4.80}
\end{align*}
$$

In Eq. (4.80), $E_{\nu}$ is to be expressed in MeV and masses $m_{i}$ and $m_{j}$ are to be given in eV .

## Oscillations between two flavours

For the sake of simplicity, let us restrict our considerations to electron neutrino ( $\nu_{e}$ ) and muon neutrino $\left(\nu_{\mu}\right)$ and the oscillations between these two flavours only. Since the mass eigenstates and the weak interaction eigenstates are connected by unitary transformation, we can consider the
unitary transformation as a rotation through an angle $\theta$ in two dimensional space xy about z axis. Accordingly, we have

$$
\left[\begin{array}{c}
\nu_{e}  \tag{4.81}\\
\nu_{\mu}
\end{array}\right]=\left[\begin{array}{rr}
\cos \theta & \sin \theta \\
-\sin \theta & \cos \theta
\end{array}\right]\left[\begin{array}{l}
\nu_{1} \\
\nu_{2}
\end{array}\right]
$$

The elements of the unitary matrix are:

$$
U_{e 1}=\cos \theta, \quad U_{e 2}=\sin \theta, \quad U_{\mu 1}=-\sin \theta, \quad U_{\mu 2}=\cos \theta,
$$

where the angle $\theta$ is known as the mixing angle. Substituting these values in Eq. (4.79), we get

$$
\begin{align*}
P_{\mu e} & =-4 U_{e 1} U_{\mu 1} U_{e 2} U_{\mu 2} \sin ^{2}\left(\frac{\pi L}{\lambda_{12}}\right) \\
& =4 \cos ^{2} \theta \sin ^{2} \theta \sin ^{2}\left(\frac{\pi L}{\lambda_{12}}\right) \\
& =\sin ^{2}(2 \theta) \sin ^{2}\left(\frac{\pi L}{\lambda_{12}}\right) .  \tag{4.82}\\
P_{e e} & =1-\sin ^{2}(2 \theta) \sin ^{2}\left(\frac{\pi L}{\lambda_{12}}\right) . \tag{4.83}
\end{align*}
$$

Equation (4.82) gives the probability of conversion of $\nu_{e}$ to $\nu_{\mu}$ during its flight over a distance $L$. Since it is a sinusoidal function of the distance of travel from the source to the detector, this phenomenon is referred to as neutrino oscillations. Eq. (4.83) gives the probability of reaching the destination without changing its flavour. Substituting the expression (4.80) for $\lambda_{12}$ in Eq. (4.82) and noting that $\pi / 2.481=1.267$, we can rewrite Eq. (4.82) as

$$
\begin{equation*}
P_{\mu e}=\sin ^{2}(2 \theta) \sin ^{2}\left(1.267 \Delta m^{2} L / E_{\nu}\right), \tag{4.84}
\end{equation*}
$$

where $\Delta m^{2}=m_{2}^{2}-m_{1}^{2}$.
Let us now discuss the physics behind Eq. (4.82) in the context of the solar neutrino puzzle. The mixing angle $\theta$ is estimated to be about $30^{\circ}$.

$$
\theta=30^{\circ} . \quad \text { Therefore } \sin ^{2}(2 \theta)=\frac{3}{4} .
$$

The oscillatory term $\sin ^{2}\left(\frac{\pi L}{\lambda_{12}}\right)$ can assume any value between 0 and 1 , depending on $L, E_{\nu}$ and $\Delta m^{2}=m_{2}^{2}-m_{1}^{2}$. The neutrino energy $E_{\nu}$ is not precise but has a spread of values, depending upon the detector that
is used to detect the neutrinos. The distance of the neutrino detector from the source is not also a constant since the earth is rotating on its own axis while describing an elliptical orbit around the Sun ${ }^{15}$. Since the events take place at a very slow rate, the data has to be collected over an extended period of time. So, it is difficult to assess exactly the value of $\sin ^{2}\left(\frac{\pi L}{\lambda_{12}}\right)$. If $\sin ^{2}\left(\frac{\pi L}{\lambda_{12}}\right) \approx 8 / 9$, then

$$
P_{\mu e} \approx \frac{3}{4} \times \frac{8}{9}=\frac{2}{3} . \quad \text { or } \quad P_{e e} \approx \frac{1}{3} .
$$

This is just an approximate value which is in agreement with the observation of the solar neutrino experiments that only one third of the solar neutrinos ( $\nu_{e}$ ) emitted by the Sun retains the same flavour, the rest getting transformed into muon neutrinos.

In the case of high energy atmospheric neutrinos, the results of SuperKamiokande experiment can be similarly interpreted by assuming the neutrino oscillation from $\nu_{\mu}$ to $\nu_{\tau}$. The high energy electron neutrinos do not change its flavour, since its oscillation length is far greater than the diameter of the earth. So, the number of high energy electron neutrinos that come from above is equal to the number of those that come from below after travelling through the earth. (vide Problems 4.5 and 4.6.)

## Oscillations between three flavours

The above restricted analysis of neutrino oscillations between two flavours is quite sufficient to explain the observations made on the solar neutrinos in Sudbury neutrino observatory and on the atmospheric neutrinos made in the Super-Kamiokande Collaboration. For the study of solar neutrinos, one can consider the oscillation of electron neutrinos into muon neutrinos. For the study of atmospheric neutrinos, one may consider the oscillation of $\nu_{\mu}$ to $\nu_{\tau}$. One can gain a preliminary knowledge by restricting the study to neutrino oscillations between two flavours. But for a detailed study, one has to consider neutrino oscillations between all the three flavours ${ }^{16}$.

The unitary transformation matrix between three flavours is given by Pontecorvo-Maki-Nakagawa-Sakata (PMNS) and it is known as the

[^33]neutrino mixing matrix $U_{P M N S}$.
\[

$$
\begin{equation*}
U_{P M N S}=R_{23} R_{13} R_{12}, \tag{4.85}
\end{equation*}
$$

\]

where $R_{23}, R_{13}, R_{12}$ denote the three rotations about the axis 1,2 and 3 respectively.

$$
\begin{align*}
& R_{23}=\left[\begin{array}{rrr}
1 & 0 & 0 \\
0 & c_{23} & s_{23} \\
0 & -s_{23} & c_{23}
\end{array}\right], R_{13}=\left[\begin{array}{ccc}
c_{13} & 0 & s_{13} e^{-i \delta} \\
0 & 1 & 0 \\
-s_{13} e^{i \delta} & 0 & c_{13}
\end{array}\right], \\
& R_{12}=\left[\begin{array}{rrr}
c_{12} & s_{12} & 0 \\
-s_{12} & c_{12} & 0 \\
0 & 0 & 1
\end{array}\right], \tag{4.86}
\end{align*}
$$

where

$$
\begin{equation*}
c_{i j}=\cos \theta_{i j}, \quad s_{i j}=\sin \theta_{i j} \tag{4.87}
\end{equation*}
$$

The angles $\theta_{12}, \theta_{13}, \theta_{23}$ are the three mixing angles and $\delta$ is the CP violating phase ${ }^{17}$. Thus we get

$$
\begin{align*}
& U_{P M N S}= \\
& \qquad\left[\begin{array}{ccc}
c_{12} c_{13} & c_{13} s_{12} & s_{13} e^{-i \delta} \\
-c_{23} s_{12}-c_{12} s_{13} s_{23} e^{i \delta} & c_{12} c_{23}-s_{12} s_{13} s_{23} e^{i \delta} & c_{13} s_{23} \\
s_{12} s_{23}-c_{12} c_{23} s_{13} e^{i \delta} & -c_{23} s_{12} s_{13} e^{i \delta}-c_{12} s_{23} & c_{13} c_{23}
\end{array}\right] \tag{4.88}
\end{align*}
$$

From a large number of experiments, so far conducted, the following values for the mixing angles and the mass-square differences have been obtained:

$$
\begin{aligned}
& \theta_{12}=35^{\circ}, \quad \theta_{23}=45^{\circ}, \quad \theta_{13}<10^{\circ} . \\
& \Delta_{12}^{2}=7.6 \times 10^{-5} \mathrm{eV}^{2}, \quad \Delta_{23}^{2}=2.4 \times 10^{-3} \mathrm{eV}^{2}, \quad \delta_{C P}=0 .
\end{aligned}
$$

Since the neutrino oscillations depend only on the mass square differences, it is not possible to obtain their absolute values. One can choose the normal hierarchy

$$
m_{3}>m_{2}>m_{1}
$$

or the inverted hierarchy

$$
m_{2}>m_{1}>m_{3}
$$

for the masses. At present, there is no compelling reason to favour one hierarchy or the other.

[^34]
### 4.6 Neutrino tomography

Neutrino oscillations occur due to a difference in phase that occurs between the wave-packets, representing the mass eigenstates. This phase difference occurs because each wave-packet propagates with different velocity due to the mass differences. This is what happens in the vacuum. In matter, however, the phase difference is determined by the total energy of the mass eigenstate. If the neutrino is propagating in a potential, $V$, then the total energy of a state is $E+V$. If the potential is different for different neutrino flavours, then the phase difference depends also on the interaction potential. Thus the neutrino oscillation depends on the matter through which it propagates. This is known as the Mikheyev-SmirnovWolfenstein (MSW) effect.

Since neutrinos can travel huge distances without being absorbed, they can be used to probe the interior of earth and the density profile of other bulk matters by observing the flavour composition of the emerging beam of neutrinos. This opens up the possibility of developing a powerful method of probe known as Neutrino Tomography.

### 4.7 India-based neutrino observatory

India has been one of the pioneers in neutrino physics. The very first detection of atmospheric neutrinos was made in the Kolar Gold Field (KGF) mines in South India in 1965. These are the neutrinos produced in the upper atmosphere by cosmic rays and hence are called atmospheric neutrinos. The KGF laboratory was closed in 1992 because of the closure of the KGF mines.

Much has happened during this period and later elsewhere - notably in USA, Canada and Japan - by setting up dedicated laboratories for neutrino research. The discovery of neutrino oscillations and their confirmation received Nobel prizes in the years 2002 and 2015. The Standard Model of elementary particles, which has been very successful in describing almost all the known phenomena hitherto, envisages massless neutrinos but the discovery of neutrino oscillations requires to endow the neutrinos with mass, although small. This necessitates to look for a theory beyond the Standard Model and the neutrino laboratories have attracted a greater attention and it is hoped that they will play a greater role in shaping the new theory.

The India-based Neutrino Observatory (INO) has been conceived with
this objective in view. The INO Laboratory will be set up in a cavern under a rocky mountain in the Bodi West Hills region in the Theni district, about 110 Km west of Madurai in Tamil Nadu.

## Review Questions

4.1 Obtain the solutions of the Dirac equation for neutrino with zero rest mass and show that they are the eigenfunctions of the chirality and helicity operators. Using the experimental observation that the neutrino is lefthanded and the antineutrino is right-handed, how do you arrive at the $\mathrm{V}-\mathrm{A}$ form of weak interaction?
4.2 Outline Weyl's two-component theory of neutrino and show that it is equivalent to the four-component Dirac theory for the neutrino, if you choose the Weyl representation for the Dirac matrices $\alpha$ and $\beta$.
4.3 Write down the Dirac matrices $\alpha$ and $\beta$ in Dirac or standard representation and also in Weyl representation. Show that one can go from one representation to another by means of unitary transformation.
4.4 What is the condition to be fulfilled for neutrino to be a Majorana particle? Does it allow mass for the left-handed neutrino? Is there any way of checking experimentally whether neutrino is a Majorana particle?
4.5 What is meant by neutrino oscillations? Discuss briefly the various experiments conducted with solar and atmospheric neutrinos that indicate neutrino oscillations.
4.6 Discuss the theory of neutrino oscillations and obtain an expression for neutrino oscillations.

## Problems

4.1 Illustrate the neutrinoless double beta decay by means of Feynman diagrams, treating nucleon (a) as an elementary particle and (b) as a composite particle consisting of quarks.
4.2 In the Japanese Kamiokande experiment, the interaction of solar neutrinos with electrons in the detector consisting of pure water in a tank is studied. Draw the Feynman diagrams for the interaction of electron neutrino and also for the converted muon neutrino with electron. Show that the electron neutrino can undergo both charged and neutral current weak interactions whereas the muon interaction can have only neutral current weak interaction.
4.3 In the Sudbury Neutrino Observatory in Canada, the interaction of solar neutrinos with deuterium in the heavy water is studied via the following reactions:

$$
\begin{aligned}
& \nu_{e}+d \rightarrow p+p+e^{-} . \\
& \nu_{x}+d \rightarrow p+n+\nu_{x},
\end{aligned} \quad x=e, \mu, \tau
$$

Draw the Feynman diagrams for these weak interactions.
4.4 It is known that the solar neutrino oscillates from electron neutrino to muon neutrino as it travels from the Sun to the earth. Find the oscillation length for the solar neutrino of energy 8 MeV , given that the mass-squared difference is

$$
\Delta m^{2}=m_{\nu_{\mu}}^{2}-m_{\nu_{e}}^{2}=7.1 \times 10^{-5} \mathrm{eV}^{2}
$$

4.5 It is known that the atmospheric neutrinos consist of high energy muon neutrinos and electron neutrinos in the GeV energy range. Find the oscillation lengths of these neutrinos, given that $m_{\nu_{\mu}}^{2}-m_{\nu_{e}}^{2}=7.1 \times 10^{-5} \mathrm{eV}^{2}$ and $\left|m_{\nu_{\tau}}^{2}-m_{\nu_{\mu}}^{2}\right|=2.35 \times 10^{-3} \mathrm{eV}^{2}$.
4.6 In the Super-Kamiokande experiment, it is found that the number of electron neutrinos that come from below is almost the same as those that come from above, whereas the number of muon neutrinos that come from below is much less than that come from above. The neutrinos that come from below travel through the earth before reaching the detector and the earth's diameter is approximately $13,000 \mathrm{Km}$. Explain this phenomenon.
4.7 Given the following values for the mixing angles:

$$
\theta_{12}=30^{\circ}, \quad \theta_{23}=45^{\circ}, \quad \theta_{13}=0^{\circ}
$$

construct the $3 \times 3$ unitary matrix $U_{P M N S}$, assuming $\delta_{C P}=0$. Express the mass eigenstates of the neutrinos in terms of their flavour eigenstates.

## Solutions to Problems

4.1 Let us consider the neutrinoless double-beta decay ( $0 \nu \beta \beta$-decay)

$$
{ }^{76} \mathrm{Ge}_{32} \rightarrow{ }^{76} \mathrm{Se}_{32}+2 e^{-},
$$

in which two neutrons are converted into two protons with the emission of two electrons. This can be illustrated by the following Feynman diagrams.

Two neutrons in the nucleus are converted into protons, each by the emission of an intermediate vector boson $W^{-}$. One $W^{-}$emits an electron and a neutrino. The emitted neutrino is absorbed by the other $W^{-}$and forms an electron in the final state. This is illustrated in Fig. 4.3 (a).


Figure 4.3: Feynman diagrams representing neutrinoless double beta decay, treating (a) nucleons as elementary particles (b) nucleons as composite particles consisting of quarks.

If one treats the nucleon as a composite state of quarks, then the neutron consists of two d quarks and one u quark whereas the proton consists of two u quarks and one d quark. In each neutron, one d quark is converted into a u quark by weak interaction. This is illustrated in Fig. 4.3 (b).
4.2 In the case of weak interaction

$$
\nu_{e}+e^{-} \rightarrow \nu_{e}+e^{-}
$$

both charged current and neutral current weak interactions are possible. The conversion of $\nu_{e}$ to $e^{-}$is possible by the emission of charged intermediate vector boson $W^{+}$. The relevant Feynman diagrams are


Figure 4.4: $\nu_{e}-e^{-}$scattering through the charged current and neutral current weak interactions by exchange of $W^{+}$and $Z^{0}$ intermediate vector bosons.

In the case of muon neutrino, $\nu_{\mu}$ cannot change to $e^{-}$due to flavour conservation. So, only the neutral current weak interaction is possible between $\nu_{\mu}$ and $e^{-}$and the relevant Feynman diagram is given below:


Figure 4.5: $\nu_{\mu}-e^{-}$scattering through the neutral current weak interaction by exchange of $Z^{0}$ intermediate vector boson.

### 4.3 The reaction

$$
\nu_{e}+d \rightarrow p+p+e^{-}
$$

corresponds to a charged current weak interaction, since it involves the exchange of $W^{+}$intermediate vector boson. The reaction

$$
\nu_{x}+d \rightarrow p+n+\nu_{x}, \quad x=e, \mu, \tau
$$

involves only the neutral current intermediate vector boson $Z^{0}$. This interaction is independent of the neutrino flavour.


Figure 4.6: (a) The reaction $\nu_{e}+d \rightarrow e^{-}+p+p$ through the exchange of $W^{+}$ intermediate vector boson (b) The reaction $\nu_{x}+d \rightarrow \nu_{x}+p+n$ through the exchange of $Z^{0}$ intermediate vector boson.
4.4 It has been shown that the probability of oscillation from electron neutrino to muon neutrino can be represented by the formula (4.83).

$$
P_{\mu e}=\sin ^{2} 2 \theta \sin ^{2}\left(\frac{\pi L}{\lambda_{21}}\right)
$$

with

$$
\lambda_{21}=2.481\left(\frac{E_{\nu}}{\mathrm{MeV}}\right)\left(\frac{(\mathrm{eV})^{2}}{m_{1}^{2}-m_{2}^{2}}\right)
$$

The $\sin ^{2}\left(\frac{\pi L}{\lambda_{21}}\right)$ is a sinusoidal function of $L$, the distance travelled. Its value oscillates between 0 and 1 as $L$ increases.

$$
\sin ^{2}\left(\frac{\pi L}{\lambda_{21}}\right)=0 \rightarrow 1 \rightarrow 0 \rightarrow 1 \cdots, \text { as } L \text { increases. }
$$

Consequently, the neutrino oscillates between the two flavours as it travels.

$$
\begin{aligned}
P_{\mu e} & =0, & & \text { if } \quad L=0 \\
& =\sin ^{2} 2 \theta, & & \text { if } \frac{\pi L}{\lambda_{21}}=\frac{\pi}{2} \\
& =0, & & \text { if } \frac{\pi L}{\lambda_{21}}=\pi
\end{aligned}
$$

The oscillation length $L_{\text {osc }}$ is given by

$$
\begin{aligned}
\frac{\pi L_{\mathrm{osc}}}{\lambda_{21}} & =\pi \\
L_{\mathrm{osc}}=\lambda_{21} & =\frac{2.481 \times 8}{7.1 \times 10^{-5}} \mathrm{~m} \\
& =279.55 \mathrm{Km}
\end{aligned}
$$

4.5 Restricting oneself to the two-flavour oscillation, let us estimate the oscillation lengths of 1 GeV muon neutrino and electron neutrino.

$$
\begin{aligned}
L_{\mathrm{osc}}(\mu \rightarrow \tau) & =\frac{2.481 \times 10^{3}}{2.35 \times 10^{-3}} \mathrm{~m} \\
& =1055.7 \mathrm{Km} \\
L_{\mathrm{osc}}(e \rightarrow \mu) & =\frac{2.481 \times 10^{3}}{7.1 \times 10^{-5}} \mathrm{~m} \\
& =3.4944 \times 10^{4} \mathrm{Km}
\end{aligned}
$$

4.6 The atmospheric neutrinos are of energy in the GeV range. For the one GeV electron neutrino, the oscillation length is approximately $3.5 \times 10^{4}$ Km , whereas for one GeV muon neutrino, the oscillation length is approximately 1000 Km . The oscillation length for the electron neutrino is much larger than the earth's diameter and so it passes through the earth without any change of flavour. In the case of muon neutrino, the oscillation length is approximately 1000 Km which is much smaller than the earth' diameter of $13,000 \mathrm{Km}$ and so it undergoes change of flavour to tau neutrino as it passes through the earth.
4.7 The PMNS unitary matrix is given by Eq. (4.88) whose elements can be determined from the given data:

$$
\begin{array}{ll}
s_{12}=\sin \theta_{12}=\sin 30^{\circ}=\frac{1}{2} & c_{12}=\cos \theta_{12}=\cos 30^{\circ}=\frac{\sqrt{3}}{2} \\
s_{23}=\sin \theta_{23}=\sin 45^{\circ}=\frac{1}{\sqrt{2}} & c_{23}=\cos \theta_{23}=\cos 45^{\circ}=\frac{1}{\sqrt{2}} \\
s_{13}=\sin \theta_{13}=\sin 0^{\circ}=0 & c_{13}=\cos \theta_{13}=\cos 0^{\circ}=1
\end{array}
$$

Substituting these values in the U-matrix (4.88), we get

$$
U=\sqrt{\frac{1}{8}}\left[\begin{array}{rrr}
\sqrt{6} & \sqrt{2} & 0 \\
-1 & \sqrt{3} & 2 \\
1 & -\sqrt{3} & 2
\end{array}\right]
$$

and

$$
U^{\dagger}=\sqrt{\frac{1}{8}}\left[\begin{array}{rrr}
\sqrt{6} & -1 & 1 \\
\sqrt{2} & \sqrt{3} & -\sqrt{3} \\
0 & 2 & 2
\end{array}\right]
$$

The mass eigenstates are given by

$$
\left[\begin{array}{l}
\nu_{1} \\
\nu_{2} \\
\nu_{3}
\end{array}\right]=U^{\dagger}\left[\begin{array}{c}
\nu_{e} \\
\nu_{\mu} \\
\nu_{\tau}
\end{array}\right]
$$

Explicitly, the mass eigenstates of the neutrinos can be written in terms of their flavour eigenstates as given below:

$$
\begin{aligned}
\nu_{1} & =\sqrt{\frac{3}{4}} \nu_{e}-\sqrt{\frac{1}{8}}\left(\nu_{\mu}-\nu_{\tau}\right) \\
\nu_{2} & =\sqrt{\frac{1}{4}} \nu_{e}+\sqrt{\frac{3}{8}}\left(\nu_{\mu}-\nu_{\tau}\right) \\
\nu_{3} & =\sqrt{\frac{1}{2}}\left(\nu_{\mu}+\nu_{\tau}\right)
\end{aligned}
$$

## Chapter 5

## The Propagation Kernel and Feynman Diagrams

In this chapter, we shall introduce the concept of propagation kernel and discuss how Feynman modified it in such a way that the positive energy states in relativistic Dirac theory are propagated forward in time whereas the negative energy states are propagated backwards in time. This new view-point has completely revolutionized the study of quantum electrodynamics and enabled the visualization of quantum electrodynamic processes by means of Feynman diagrams. It is said that more than Feynman, it was Freeman Dyson who has popularized the Feynman Diagrams and showed that Feynman's approach is equivalent to the field theoretical method.

For the study of a given electrodynamic process, one has to draw all possible Feynman diagrams for the process and calculate the transition probability by formulating a set of rules known as Feynman rules.

### 5.1 The Propagation Kernel

### 5.1.1 In non-relativistic Schrödinger theory

First let us consider the non-relativistic Schrödinger equation

$$
\begin{equation*}
i \hbar \frac{\partial \Psi(\boldsymbol{x}, t)}{\partial t}=H \Psi(\boldsymbol{x}, t), \tag{5.1}
\end{equation*}
$$

and introduce the concept of propagation kernel. It is an extension of the Green's function technique ${ }^{1}$, introduced in the theory of scattering. Equation (5.1) gives the infinitesimal change in the wave function $\Psi$ due to a small increment in time. To obtain $\Psi\left(\boldsymbol{x}_{2}, t_{2}\right)$ from $\Psi\left(\boldsymbol{x}_{1}, t_{1}\right)$, one has to go through a series of successive incremental steps but Feynman chose to take a great leap forward with the help of a propagation kernel. The propagation kernel $K\left(\boldsymbol{x}_{2} t_{2} ; \boldsymbol{x}_{1} t_{1}\right)$ describes the evolution of the state function $\Psi(\boldsymbol{x}, t)$ and it is defined by

$$
\begin{equation*}
\Psi\left(\boldsymbol{x}_{2}, t_{2}\right)=\int d^{3} x_{1} K\left(\boldsymbol{x}_{2} t_{2}: \boldsymbol{x}_{1} t_{1}\right) \Psi\left(\boldsymbol{x}_{1}, t_{1}\right), \quad t_{2} \geq t_{1} . \tag{5.2}
\end{equation*}
$$

First let us assume that $H$ is a time independent operator with eigenfunctions $\psi_{n}(\boldsymbol{x})$ and eigenvalues $E_{n}$, such that

$$
\begin{equation*}
H \psi_{n}(\boldsymbol{x})=E_{n} \psi_{n}(\boldsymbol{x}) \tag{5.3}
\end{equation*}
$$

The eigenfunctions are orthonormal and obey the closure property.

$$
\begin{align*}
\int d^{3} \boldsymbol{x} \psi_{n}^{*}(\boldsymbol{x}) \psi_{m}(\boldsymbol{x}) & =\delta_{m n}  \tag{5.4}\\
\sum_{n} \psi_{n}(\boldsymbol{x}) \psi_{n}\left(\boldsymbol{x}^{\prime}\right) & =\delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right) \tag{5.5}
\end{align*}
$$

We can now expand $\Psi\left(\boldsymbol{x}_{1}, t_{1}\right)$ in terms of the complete set of orthonormal eigenfunctions $\psi_{n}(\boldsymbol{x})$.

$$
\begin{equation*}
\Psi\left(\boldsymbol{x}_{1}, t_{1}\right)=\sum_{n} c_{n} \psi_{n}\left(\boldsymbol{x}_{1}\right) e^{-i E_{n} t_{1} / \hbar} \tag{5.6}
\end{equation*}
$$

From Eq. (5.6), the expansion coefficients $c_{n}$ can be obtained by using the orthonormal and closure properties of $\psi_{n}(\boldsymbol{x})$.

$$
\begin{equation*}
c_{n}=\int \psi_{n}^{*}\left(\boldsymbol{x}_{1}\right) \Psi\left(\boldsymbol{x}_{1}, t_{1}\right) e^{i E_{n} t_{1} / \hbar} d^{3} x_{1} . \tag{5.7}
\end{equation*}
$$

Using these expansion coefficients, the state vector $\Psi\left(\boldsymbol{x}_{2}, t_{2}\right)$ can be written as

$$
\begin{align*}
\Psi\left(\boldsymbol{x}_{2}, t_{2}\right) & =\sum_{n} c_{n} \psi_{n}\left(\boldsymbol{x}_{2}\right) e^{-i E_{n} t_{2} / \hbar}  \tag{5.8}\\
& =\sum_{n} \int \psi_{n}^{*}\left(\boldsymbol{x}_{1}\right) \Psi\left(\boldsymbol{x}_{1}, t_{1}\right) \psi_{n}\left(\boldsymbol{x}_{2}\right) e^{-i E_{n}\left(t_{2}-t_{1}\right) / \hbar} d^{3} \boldsymbol{x}_{1} . \tag{5.9}
\end{align*}
$$

[^35]Comparing Eqs. (5.2) and (5.9), we obtain an expression for the propagation kernel $K\left(\boldsymbol{x}_{2} t_{2} ; \boldsymbol{x}_{1} t_{1}\right)$.

$$
\begin{align*}
K\left(\boldsymbol{x}_{2} t_{2} ; \boldsymbol{x}_{1} t_{1}\right) & =\sum_{n} \psi_{n}\left(\boldsymbol{x}_{2}\right) \psi_{n}^{*}\left(\boldsymbol{x}_{1}\right) e^{-i E_{n}\left(t_{2}-t_{1}\right) / \hbar}, \quad t_{2}>t_{1} \\
& =\sum_{n} \chi_{n}\left(\boldsymbol{x}_{2}, t_{2}\right) \chi_{n}^{*}\left(\boldsymbol{x}_{1}, t_{1}\right) \theta\left(t_{2}-t_{1}\right), \tag{5.10}
\end{align*}
$$

where

$$
\begin{align*}
\chi_{n}(\boldsymbol{x}, t) & =\psi_{n}(\boldsymbol{x}) e^{-i E_{n} t / \hbar},  \tag{5.11}\\
\theta\left(t_{2}-t_{1}\right) & =\left\{\begin{array}{ll}
1, & \text { if } t_{2}>t_{1} \\
0, & \text { if } t_{2}<t_{1}
\end{array} .\right. \tag{5.12}
\end{align*}
$$

It is now possible to obtain a differential equation which the propagation kernel obeys. Since $i \hbar \frac{\partial}{\partial t_{2}} \chi_{n}\left(\boldsymbol{x}_{2}, t_{2}\right)=H \chi_{n}\left(\boldsymbol{x}_{2}, t_{2}\right)$, we obtain

$$
\begin{align*}
& \left\{i \hbar \frac{\partial}{\partial t_{2}}-H\left(\boldsymbol{x}_{2}\right)\right\} K\left(\boldsymbol{x}_{2} t_{2} ; \boldsymbol{x}_{1} t_{1}\right) \\
& \quad=i \hbar \sum_{n} \chi_{n}\left(\boldsymbol{x}_{2}, t_{2}\right) \chi_{n}^{*}\left(\boldsymbol{x}_{1}, t_{1}\right) \frac{\partial}{\partial t_{2}} \theta\left(t_{2}-t_{1}\right) \\
& =i \hbar \sum_{n} \psi_{n}\left(\boldsymbol{x}_{2}\right) \psi_{n}^{*}\left(\boldsymbol{x}_{1}\right) e^{-i E_{n}\left(t_{2}-t_{1}\right) / \hbar} \delta\left(t_{2}-t_{1}\right) \\
& \quad=i \hbar \sum_{n} \psi_{n}\left(\boldsymbol{x}_{2}\right) \psi_{n}^{*}\left(\boldsymbol{x}_{1}\right) \delta\left(t_{2}-t_{1}\right) \\
& \quad=i \hbar \delta\left(\boldsymbol{x}_{2}-\boldsymbol{x}_{1}\right) \delta\left(t_{2}-t_{1}\right) . \tag{5.13}
\end{align*}
$$

The last step is obtained using the closure property (5.5). Although Eq. (5.13) is obtained assuming $H$ to be time independent, it is true in the general case ${ }^{2}$ when $H$ is also time dependent.

$$
\begin{equation*}
\left\{i \hbar \frac{\partial}{\partial t_{2}}-H\left(\boldsymbol{x}_{2}, t_{2}\right)\right\} K\left(\boldsymbol{x}_{2} t_{2} ; \boldsymbol{x}_{1} t_{1}\right)=i \hbar \delta\left(\boldsymbol{x}_{2}-\boldsymbol{x}_{1}\right) \delta\left(t_{2}-t_{1}\right) \tag{5.14}
\end{equation*}
$$

Equation (5.14) is the differential equation for the propagation kernel $K\left(\boldsymbol{x}_{2} t_{2} ; \boldsymbol{x}_{1} t_{1}\right)$ and it is equivalent to the integral form (5.2).

Now, we are in a position to develop a perturbation theory. Suppose

$$
H=H_{0}+V,
$$

[^36]where $V$ is a small perturbation. We know the solutions of the unperturbed Hamiltonian $H_{0}$. Then the propagation kernel $K_{0}$ for $V=0$ corresponds to the free particle. This enables us to obtain an integral equation for $K$ which is equivalent to the differential equation (5.14) with the boundary condition ${ }^{3}$
$$
K(\mathbf{2}, \mathbf{1})=0, \quad t_{2}<t_{1} .
$$
and develop a perturbation theory for $K(\mathbf{2}, \mathbf{1})$ in terms of $K_{0}(\mathbf{2}, \mathbf{1})$. The integral equation for $K(\mathbf{2}, \mathbf{1})$ is given below:
\[

$$
\begin{equation*}
K(\mathbf{2}, \mathbf{1})=K_{0}(\mathbf{2}, \mathbf{1})-\frac{i}{\hbar c} \int K_{0}(\mathbf{2}, \mathbf{3}) V(\mathbf{3}) K(\mathbf{3}, \mathbf{1}) d^{4} x_{3}, \tag{5.15}
\end{equation*}
$$

\]

where $d^{4} x$ is the element of the four dimensional space-time volume ${ }^{4}$.

$$
\begin{equation*}
d^{4} x=c d t d x d y d z=d x_{0} d x_{1} d x_{2} d x_{3} . \tag{5.16}
\end{equation*}
$$

It can be easily verified that the integral equation (5.15) is equivalent to the differential equation

$$
\begin{equation*}
\left\{i \hbar \frac{\partial}{\partial t_{2}}-H_{0}(\mathbf{2})\right\} K(\mathbf{2}, \mathbf{1})=i \hbar \delta^{4}(\mathbf{2}, \mathbf{1})+V(\mathbf{2}) K(\mathbf{2}, \mathbf{1}) . \tag{5.17}
\end{equation*}
$$

Equation (5.17) is identical with Eq. (5.14). The integral equation (5.15) can be written as a perturbation series by repeated application of the expansion for $K$ occurring on the right hand side of Eq. (5.15).

$$
\begin{aligned}
K(\mathbf{2}, \mathbf{1})= & K_{0}(\mathbf{2}, \mathbf{1})+\left(\frac{-i}{\hbar c}\right) \int K_{0}(\mathbf{2}, \mathbf{3}) V(\mathbf{3}) K_{0}(\mathbf{3}, \mathbf{1}) d^{4} x_{3} \\
& +\left(\frac{-i}{\hbar c}\right)^{2} \int K_{0}(\mathbf{2}, \mathbf{3}) V(\mathbf{3}) K_{0}(\mathbf{3}, \mathbf{4}) V(\mathbf{4}) K_{0}(\mathbf{4}, \mathbf{1}) d^{4} x_{3} d^{4} x_{4}
\end{aligned}
$$

$$
\begin{equation*}
+\cdots \tag{5.18}
\end{equation*}
$$

### 5.1.2 In relativistic Dirac theory

Consider the time-dependent Dirac equation in Feynman notation for a free particle ${ }^{5}$ using natural units ( $\hbar=c=1$ ).

$$
\begin{equation*}
(i \not \nabla-m) \Psi(\boldsymbol{x}, t)=0 . \tag{5.19}
\end{equation*}
$$

[^37]Here $\Psi\left(\boldsymbol{x}_{1}, t_{1}\right)$ and $\Psi\left(\boldsymbol{x}_{2}, t_{2}\right)$ are four component wave functions and the propagation kernel $K\left(\boldsymbol{x}_{2}, t_{2} ; \boldsymbol{x}_{1}, t_{1}\right)$ is a $(4 \times 4)$ matrix. The differential equation that the propagation kernel obeys is analogous to Eq. (5.14).

$$
\begin{equation*}
(i \nabla-m) K_{0}(\mathbf{2}, \mathbf{1})=i \delta^{4}(\mathbf{2}, \mathbf{1}) . \tag{5.20}
\end{equation*}
$$

Expanding $\Psi$ in terms of the complete set of eigenfunctions of (5.19), which correspond to both positive and negative energy eigenvalues, we obtain

$$
K_{0}(\mathbf{2}, \mathbf{1})= \begin{cases}\sum_{E_{n}>0} \psi_{n}\left(\boldsymbol{x}_{2}\right) \bar{\psi}_{n}\left(\boldsymbol{x}_{1}\right) e^{-i E_{n}\left(t_{2}-t_{1}\right)} &  \tag{5.21}\\ +\sum_{E_{n}<0} \psi_{n}\left(\boldsymbol{x}_{2}\right) \bar{\psi}_{n}\left(\boldsymbol{x}_{1}\right) e^{-i E_{n}\left(t_{2}-t_{1}\right)}, & t_{2}>t_{1}, \\ 0, & t_{2}<t_{1}\end{cases}
$$

For relativistic convenience, one deals with adjoint wave functions $\bar{\psi}$ instead of $\psi^{\dagger}$ in Eq. (5.21). The solution (5.21) is unsatisfactory since it corresponds to the one-electron theory instead of the correct hole theory. It admits the transition of the electron from the positive energy state to the negative energy state for $t_{2}>t_{1}$. This is clearly unphysical since, according to the hole theory, all the negative energy states are full and hence not available for the electron after scattering. So, for $t_{2}>t_{1}, K_{0}$ should contain only positive energy states. Each term $\psi_{n}\left(\boldsymbol{x}_{2}\right) \bar{\psi}_{n}\left(\boldsymbol{x}_{1}\right) e^{-i E_{n}\left(t_{2}-t_{1}\right)}$ is a solution of the homogeneous equation if it is used for all times $t_{2}$. So, it is permitted to subtract the sum over all the negative energy states for both $t_{2}>t_{1}$ and $t_{2}<t_{1}$. The resulting propagation kernel $K_{F}$, known as Feynman's Propagation Kernel, is given by

$$
K_{F}(\mathbf{2}, \mathbf{1})= \begin{cases}\sum_{E_{n}>0} \psi_{n}\left(\boldsymbol{x}_{2}\right) \bar{\psi}_{n}\left(\boldsymbol{x}_{1}\right) e^{-i E_{n}\left(t_{2}-t_{1}\right)}, & t_{2}>t_{1}  \tag{5.22}\\ -\sum_{E_{n}<0} \psi_{n}\left(\boldsymbol{x}_{2}\right) \bar{\psi}_{n}\left(\boldsymbol{x}_{1}\right) e^{-i E_{n}\left(t_{2}-t_{1}\right)}, & t_{2}<t_{1}\end{cases}
$$

Please note that $e^{-i E_{n}\left(t_{2}-t_{1}\right)}=e^{-i\left|E_{n}\right|\left|t_{2}-t_{1}\right|}$, for both $t_{2}>t_{1}$ and $t_{2}<t_{1}$. The perturbation expansion (5.18) is still valid if $K_{0}$ is replaced by $K_{F}$. Only $V$ is now a $4 \times 4$ matrix.

$$
\begin{aligned}
& \overline{\text { since } p=\gamma_{\mu} p_{\mu}=\gamma_{0} E-\boldsymbol{\gamma} \cdot \boldsymbol{p} \text {, with } E=i \frac{\partial}{\partial t}=i \frac{\partial}{\partial x_{0}} \text { and } \boldsymbol{p}=-i \boldsymbol{\nabla} \text {. Thus }} \\
& p=i\left\{\gamma_{0} \frac{\partial}{\partial x_{0}}+\gamma_{x} \frac{\partial}{\partial x}+\gamma_{y} \frac{\partial}{\partial y}+\gamma_{z} \frac{\partial}{\partial z}\right\}=i \nabla \boldsymbol{\nabla} \text {, where } \not \nabla=\gamma_{0} \frac{\partial}{\partial x_{0}}+\gamma_{x} \frac{\partial}{\partial x}+\gamma_{y} \frac{\partial}{\partial y}+\gamma_{z} \frac{\partial}{\partial z} .
\end{aligned}
$$

The choice of the propagation kernel (5.22) causes the positive energy components to be propagated forward in time while the negative energy components are propagated backwards in time. This is the most notable contribution of Feynman, which has resulted in the birth of Feynman's positron theory, according to which the electrons are propagated forward in time whereas the positrons correspond to negative energy electrons that are propagated backwards in time.

### 5.1.3 In momentum representation

The Feynman propagation kernel $K_{F}$ obeys the differential equation ${ }^{6}$

$$
\begin{equation*}
(i \not \nabla-m) K_{F}(\mathbf{2}, \mathbf{1})=i \delta^{4}(\mathbf{2}, \mathbf{1}) \tag{5.23}
\end{equation*}
$$

If $S_{+}(\mathbf{p})$ is the Fourier transform of $K_{F}$, then

$$
\begin{equation*}
K_{F}(\mathbf{2}, \mathbf{1})=\int_{-\infty}^{+\infty} S_{+}(\mathbf{p}) e^{-i \mathbf{p} \cdot\left(\mathbf{x}_{2}-\mathbf{x}_{1}\right)} d^{4} p \tag{5.24}
\end{equation*}
$$

where

$$
d^{4} p=d p_{0} d p_{x} d p_{y} d p_{z}=d p_{0} d p_{1} d p_{2} d p_{3} .
$$

Denoting the four-dimensional delta function in the form of Fourier transform,

$$
\begin{equation*}
\delta^{4}(\mathbf{2}, \mathbf{1})=\frac{1}{(2 \pi)^{4}} \int_{-\infty}^{+\infty} e^{-i \mathbf{p} \cdot\left(\mathbf{x}_{2}-\mathbf{x}_{1}\right)} d^{4} p \tag{5.25}
\end{equation*}
$$

Eq. (5.23) can be written as

$$
\begin{equation*}
(i \not \nabla-m) \int_{-\infty}^{+\infty} S_{+}(\mathbf{p}) e^{-i \mathbf{p} \cdot\left(\mathbf{x}_{2}-\mathbf{x}_{1}\right)} d^{4} p=\frac{i}{(2 \pi)^{4}} \int_{-\infty}^{+\infty} e^{-i \mathbf{p} \cdot\left(\mathbf{x}_{2}-\mathbf{x}_{1}\right)} d^{4} p \tag{5.26}
\end{equation*}
$$

Equating the integrands in Eq. (5.26), we obtain

$$
\begin{equation*}
(\not p-m) S_{+}(\mathbf{p})=\frac{i}{(2 \pi)^{4}} \tag{5.27}
\end{equation*}
$$

This yields the propagation kernel in the momentum representation.

$$
\begin{equation*}
S_{+}(\mathbf{p})=\frac{i}{(2 \pi)^{4}} \frac{1}{(\not p-m)}=\frac{i}{(2 \pi)^{4}} \frac{(\not p+m)}{\mathbf{p}^{2}-m^{2}}, \tag{5.28}
\end{equation*}
$$

[^38]since
$$
(\not p+m)(\not p-m)=\not p^{2}-m^{2} \quad \text { and } \quad \not p^{2}=p_{0}^{2}-\boldsymbol{p}^{2}=\mathbf{p}^{2} .
$$

The foregoing discussion yields an expression for the propagation kernel $K_{F}(2,1)$.

$$
\begin{equation*}
K_{F}(\mathbf{2}, \mathbf{1})=\frac{i}{(2 \pi)^{4}} \int \frac{1}{(\not p-m)} e^{-i \mathbf{p} \cdot\left(\mathbf{x}_{2}-\mathbf{x}_{1}\right)} d^{4} p \tag{5.29}
\end{equation*}
$$

Equation (5.29) is obtained from Eqs. (5.24) and (5.28).

### 5.1.4 Interaction with electromagnetic field

Let us consider the interaction of an electron with an external electromagnetic field which is represented by an interaction potential

$$
\begin{equation*}
V=e \gamma_{\mu} A_{\mu}=e \not A \tag{5.30}
\end{equation*}
$$

The Dirac equation for an electron in an electromagnetic field is given by

$$
\begin{equation*}
(i \not \nabla-m) \Psi(\mathbf{x})=e \neq A \Psi(\mathbf{x}) \tag{5.31}
\end{equation*}
$$

The propagation kernel $K_{F}^{A}$ for the electron in an electromagnetic field satisfies the differential equation

$$
\begin{equation*}
(i \not \nabla-e \not \subset-m) K_{F}^{A}(2,1)=i \delta^{4}(2,1) \tag{5.32}
\end{equation*}
$$

The propagation kernel has only positive energy components for $t_{2}>t_{1}$ and only negative energy components for $t_{2}<t_{1}$. It satisfies the integral equation

$$
\begin{equation*}
K_{F}^{A}(\mathbf{2}, \mathbf{1})=K_{F}(\mathbf{2}, \mathbf{1})-i e \int d^{4} x_{3} K_{F}(\mathbf{2}, \mathbf{3}) \notin(\mathbf{3}) K_{F}^{A}(\mathbf{3}, \mathbf{1}) \tag{5.33}
\end{equation*}
$$

The perturbation expansion of this integral equation is given below.

$$
\begin{align*}
K_{F}^{A}(\mathbf{2}, \mathbf{1})= & K_{F}(\mathbf{2}, \mathbf{1})-i e \int d^{4} x_{3} K_{F}(\mathbf{2}, \mathbf{3}) A(\mathbf{3}) K_{F}(\mathbf{3}, \mathbf{1}) \\
& +(-i e)^{2} \int d^{4} x_{3} \int d^{4} x_{4} K_{F}(\mathbf{2}, \mathbf{3}) A(\mathbf{3}) K_{F}(\mathbf{3}, \mathbf{4}) A(\mathbf{4}) K_{F}(\mathbf{4}, \mathbf{1}) \\
& +\cdots . \tag{5.34}
\end{align*}
$$

Let us now write down the Dirac wave function $\Psi(\mathbf{2})$ at the space-time point $\mathbf{2}$ in terms of the wave function $\Psi(\mathbf{1})$ at the space-time point $\mathbf{1}$.

$$
\begin{equation*}
\Psi(\mathbf{2})=\int K_{F}^{A}(\mathbf{2}, \mathbf{1}) \beta \Psi(\mathbf{1}) d^{3} x_{1} \tag{5.35}
\end{equation*}
$$

For the non-interacting case

$$
\begin{equation*}
\Psi_{0}(\mathbf{2})=\int K_{F}(\mathbf{2}, \mathbf{1}) \beta \Psi_{0}(\mathbf{1}) d^{3} x_{1} \tag{5.36}
\end{equation*}
$$

where $\Psi_{0}$ is a solution of the free Dirac equation.

### 5.2 The Transition Amplitude

### 5.2.1 First order matrix element

The transition amplitude for an electron to go from some positive energy state $\Psi_{i}\left(\boldsymbol{x}_{1}, t_{1}\right)$ at time $t_{1}$ to some other positive energy state $\Psi_{f}\left(\boldsymbol{x}_{2}, t_{2}\right)$ at time $t_{2}$ is given by

$$
\begin{align*}
\mathscr{M} & =\int \Psi_{f}^{\dagger}\left(\boldsymbol{x}_{2}, t_{2}\right) K_{F}\left(\boldsymbol{x}_{2} t_{2}, \boldsymbol{x}_{1} t_{1}\right) \beta \Psi_{i}\left(\boldsymbol{x}_{1}, t_{1}\right) d^{3} x_{1} d^{3} x_{2} \\
& =\int \bar{\Psi}_{f}(\mathbf{2}) \beta K_{F}(\mathbf{2}, \mathbf{1}) \beta \Psi_{i}(\mathbf{1}) d^{3} x_{1} d^{3} x_{2} \tag{5.37}
\end{align*}
$$

If an external potential acts between $t_{1}$ and $t_{2}$, then $K_{F}$ is replaced by $K_{F}^{A}$ and one can obtain the first order transition matrix element $\mathscr{M}_{1}$ using the perturbation series (5.34).

$$
\begin{equation*}
\mathscr{M}_{1}=-i e \iiint \bar{\Psi}_{f}(2) \beta K_{F}(\mathbf{2}, \mathbf{3}) A(\mathbf{3}) K_{F}(\mathbf{3}, \mathbf{1}) \beta \Psi_{i}(\mathbf{1}) d^{3} x_{1} d^{3} x_{2} d^{4} x_{3} \tag{5.38}
\end{equation*}
$$

Since

$$
\begin{align*}
& \int K_{F}(\mathbf{3}, \mathbf{1}) \beta \Psi_{i}(\mathbf{1}) d^{3} x_{1}=\Psi_{i}(\mathbf{3})  \tag{5.39}\\
& \int \bar{\Psi}_{f}(\mathbf{2}) \beta K_{F}(\mathbf{2}, \mathbf{3}) d^{3} x_{2}=\bar{\Psi}_{f}(\mathbf{3}) \tag{5.40}
\end{align*}
$$

the first order matrix element can be rewritten as

$$
\begin{equation*}
\mathscr{M}_{1}=-i e \int \bar{\Psi}_{f}(\mathbf{3}) A(\mathbf{3}) \Psi_{i}(\mathbf{3}) d^{4} x_{3} \tag{5.41}
\end{equation*}
$$

Let us choose the initial state to be the positive energy state of an electron with four-momentum $\mathbf{p}_{1}$ and the final state to be the positive energy state of the electron with four-momentum $\mathbf{p}_{2}$.

$$
\begin{align*}
\Psi_{i}(\mathbf{x}) & =u\left(\boldsymbol{p}_{1}\right) e^{-i \mathbf{p}_{1} \cdot \mathbf{x}}  \tag{5.42}\\
\bar{\Psi}_{f}(\mathbf{x}) & =\bar{u}\left(\boldsymbol{p}_{2}\right) e^{i \mathbf{p}_{2} \cdot \mathbf{x}} \tag{5.43}
\end{align*}
$$

If $a_{\mu}(\mathbf{k})$ is the Fourier transform of $A_{\mu}(\mathbf{x})$, then

$$
\begin{equation*}
A_{\mu}(\mathbf{x})=\int a_{\mu}(\mathbf{k}) e^{-i \mathbf{k} \cdot \mathbf{x}} d^{4} k \tag{5.44}
\end{equation*}
$$

Substituting the results (5.42) - (5.44) in Eq. (5.41), we obtain

$$
\begin{equation*}
\mathscr{M}_{1}=-i e \int \bar{u}\left(\boldsymbol{p}_{2}\right) d(\mathbf{k}) u\left(\boldsymbol{p}_{1}\right) e^{i\left(\mathbf{p}_{2}-\mathbf{k}-\mathbf{p}_{1}\right) \cdot \mathbf{x}} d^{4} x d^{4} k \tag{5.45}
\end{equation*}
$$

Integrating over $d^{4} x$, we get

$$
\begin{equation*}
\mathscr{M}_{1}=-i e(2 \pi)^{4} \int \bar{u}\left(\boldsymbol{p}_{2}\right) \phi(\mathbf{k}) u\left(\boldsymbol{p}_{1}\right) \delta^{4}\left(\mathbf{p}_{2}-\mathbf{k}-\mathbf{p}_{1}\right) d^{4} k \tag{5.46}
\end{equation*}
$$

where $\delta^{4}\left(\mathbf{p}_{2}-\mathbf{k}-\mathbf{p}_{1}\right)$ denotes the energy-momentum conservation at the interaction point. The integration over $d^{4} k$ is trivial and we finally obtain the first-order matrix element in a simple form.

$$
\begin{equation*}
\mathscr{M}_{1}=-i e(2 \pi)^{4} \bar{u}\left(\boldsymbol{p}_{2}\right) d(\mathbf{k}) u\left(\boldsymbol{p}_{1}\right), \quad \text { with } \quad \mathbf{k}=\mathbf{p}_{2}-\mathbf{p}_{1} . \tag{5.47}
\end{equation*}
$$

### 5.2.2 Second order matrix element

In a similar way, we can obtain the second-order matrix element for the electron transition from the initial momentum state $\boldsymbol{p}_{1}$ to the final momentum state $\boldsymbol{p}_{2}$. Adopting the same method that we have used to obtain Eq. (5.41) for the first order matrix element $\mathscr{M}_{1}$, we obtain an expression for the second order matrix element.

$$
\begin{equation*}
\mathscr{M}_{2}=(-i e)^{2} \int d^{4} x_{3} d^{4} x_{4} \bar{\Psi}_{f}(\mathbf{4}) A(\mathbf{4}) K_{F}(\mathbf{4}, \mathbf{3}) A(\mathbf{3}) \Psi_{i}(\mathbf{3}) . \tag{5.48}
\end{equation*}
$$

Using expressions (5.42) - (5.44) and (5.29), and performing the integration over $d^{4} x_{3}$ and $d^{4} x_{4}$, we can go to the momentum representation to obtain

$$
\begin{align*}
\mathscr{M}_{2}=(-i e)^{2} & \int d^{4} p d^{4} k_{1} d^{4} k_{2}(2 \pi)^{4} \delta^{4}\left(\mathbf{p}_{2}-\mathbf{k}_{2}-\mathbf{p}\right) \delta^{4}\left(\mathbf{p}-\mathbf{k}_{1}-\mathbf{p}_{1}\right) \\
& \times \bar{u}\left(\boldsymbol{p}_{2}\right) d\left(\mathbf{k}_{2}\right) \frac{1}{\not p-m} d\left(\mathbf{k}_{1}\right) u\left(\boldsymbol{p}_{1}\right) \tag{5.49}
\end{align*}
$$

The integration over $d^{4} p$ and $d^{4} k_{2}$ can be performed trivially because of the $\delta$ - functions and we are left with only one integration.

$$
\begin{equation*}
\mathscr{M}_{2}=(-i e)^{2}(2 \pi)^{4} \int d^{4} k_{1} \bar{u}\left(\boldsymbol{p}_{2}\right) d\left(\mathbf{p}_{2}-\mathbf{p}_{1}-\mathbf{k}_{1}\right) \frac{1}{\not p_{1}+\not / k_{1}-m} d\left(\mathbf{k}_{1}\right) u\left(\boldsymbol{p}_{1}\right) . \tag{5.50}
\end{equation*}
$$

Eqs. (5.47) and (5.50) have a simple physical interpretation. Eq. (5.47) describes an electron of momentum $\boldsymbol{p}_{1}$ interacting with an electromagnetic field, represented by -ied ( $\mathbf{k}$ ), picking up a four-momentum $\mathbf{k}$ and making a transition to the final state with momentum $\boldsymbol{p}_{2}$. In a similar way, Eq. (5.50) denotes an electron, that interacts with an electromagnetic field of four-momentum $\mathbf{k}_{1}$, propagates as a free particle with a kernel $\frac{1}{\not p_{1}+k_{1}-m}$ until it interacts once again with the electromagnetic field and emerges as an electron with momentum $\boldsymbol{p}_{2}$.

### 5.3 Feynman Diagrams

The propagation of electron interacting with an external field can be conveniently depicted by Feynman diagrams as shown in Fig. 5.1. (A) gives the Feynman diagrams in space-time representation and (B) gives the Feynman diagrams in momentum-energy representation. Both first order and second order transitions are considered for illustrative purpose.

The propagation of electron from space-time point (1) to space-time point (2) is usually called a world line. The world line need not always go forward in time but sometimes it can also go backwards in time.

### 5.3.1 In space-time representation

## Diagrams A(a) and A(b)

As we follow the world line from (1) to (2) in diagrams $\mathrm{A}(\mathrm{a})$ and $\mathrm{A}(\mathrm{b})$, time always increases. So in diagram $\mathrm{A}(\mathrm{a})$, the propagation kernels $K_{F}(\mathbf{3}, \mathbf{1})$ and $K_{F}(\mathbf{2}, \mathbf{3})$ consist of sums over positive energy states only as given by Eq. (5.22) for $t_{2}>t_{1}$. Similar is the situation in case of diagram A(b) where all the propagation kernels $K_{F}(\mathbf{3}, \mathbf{1}), K_{F}(\mathbf{4}, \mathbf{3})$ and $K_{F}(\mathbf{2}, \mathbf{1})$ point towards the increasing time and hence consist of sums over positive energy states only. Diagram $\mathrm{A}(\mathrm{b})$ represents the conventional double scattering.

## Diagram A(b')

As we follow the world line, we find that we go backwards in time from $\mathbf{3}$ to 4 and hence the propagation kernel $K_{F}(\mathbf{4}, \mathbf{3})$ consists of a sum of negative energy states only. Feynman interprets this as positron propagating from 4 to 3. In the conventional view-point of increasing time, diagram $\mathrm{A}\left(\mathrm{b}^{\prime}\right)$ represents the following sequence of events. An electron-positron pair is created at space-time point 4 , the positron gets annihilated with the incoming electron at space-time point $\mathbf{3}$ and the electron propagates to $\mathbf{2}$ which is observed as the final state in scattering.

(B) In momentum-energy representation

Figure 5.1: Feynman diagrams (A) in space-time representation and (B) in momentum-energy representation representing the first order scattering (a) and the second order scattering ( b and $\mathrm{b}^{\prime}$ ). The diagrams ( b ) and ( $\mathrm{b}^{\prime}$ ) are equivalent. At each vertex, energy and momentum are conserved.

It is emphasized that both diagrams $A(b)$ and $A\left(b^{\prime}\right)$ are included in
the single Eq. (5.48) since an integration is performed over all times $t_{3}$ and $t_{4}$. So, Eq. (5.48) is usually represented by a single diagram $\mathrm{A}(\mathrm{b})$ and the diagram $\mathrm{A}\left(\mathrm{b}^{\prime}\right)$ is considered identical to it, since it can be obtained from diagram $\mathrm{A}(\mathrm{b})$ when $t_{4}<t_{3}$. This way of description satisfies the Dirac hole theory. A negative energy electron is lifted from the occupied negative energy sea to the electron final state $\mathbf{2}$ whereas the hole in the negative energy sea is filled by the incoming electron $\mathbf{1}$. What is the observed effect that diagram $\mathrm{A}\left(\mathrm{b}^{\prime}\right)$ represents? An electron is scattered from state $\mathbf{1}$ to state $\mathbf{2}$ but at the same time a positive energy electron and a negative-sea electron is exchanged. The Pauli principle demands that the matrix element be given a negative sign. Indeed the Feynman propagation kernel $K_{F}(\mathbf{2}, \mathbf{1})$, given by Eq. (5.22), has a negative sign for $t_{2}<t_{1}$ without explicitly using the Pauli principle. Thus the dynamics of Feynman's propagation kernel $K_{F}$ include the negative sign as required by the Pauli principle. By this novel approach, Feynman is able to develop a theory which can describe quantum electrodynamical processes without recourse to the complications of old-fashioned hole theory with its intermediate states. Thus Feynman's approach greatly simplifies the calculational procedure and so, it has superseded all the older methods.

### 5.3.2 In momentum-energy representation

Figure 3.1(B) presents the Feynman diagrams of the afore-discussed events in momentum-energy representation. Diagram $\mathrm{B}(\mathrm{a})$ depicts the first order scattering of an electron by an external field. An electron with four-momentum $\mathbf{p}_{1}$ interacts with an external field and picks up a fourmomentum $\mathbf{k}$ and reaches the final state with four-momentum $\mathbf{p}_{2}$. Diagrams $\mathrm{B}(\mathrm{b})$ and $\mathrm{B}\left(\mathrm{b}^{\prime}\right)$ depict the second order scattering. Both the diagrams describe an electron with an initial four-momentum $\mathbf{p}_{1}$, that is scattered by an external field, picking up a four-momentum $\mathbf{k}_{1}$ and travels as a virtual particle with a propagator $\frac{1}{\not p-m}$ and once again scattered, picking up a four-momentum $\mathbf{k}_{2}$ before reaching the final state with fourmomentum $\mathbf{p}_{2}$. Although the topology of the two diagrams $\mathrm{B}(\mathrm{b})$ and $B\left(b^{\prime}\right)$ are slightly different due to the different ways of representing the intermediate state, both of them can be described mathematically by the same expression although they can be considered physically as arising from different physical processes as explained earlier in our discussion on space-time representation.

The momentum-energy representation of the events are of great importance to us since most of the calculations are done in this representation.

At each vertex, energy and momentum are conserved but at the intermediate states, the particle does not satisfy the required energy-momentum relation $E^{2}=p^{2}+m^{2}$ and hence the intermediate state particles are sometimes called virtual particles.

For calculational purpose, the momentum-energy representation is found to be more convenient and most of the quantum electrodynamic processes are studied by drawing Feynman diagrams in energy-momentum representation.

### 5.3.3 Basic vertices in QED

Feynman diagrams in Quantum Electrodynamics (QED) are made up of vertices which are points where three lines (electron, positron and photon lines) meet representing electromagnetic interaction with electron or positron. The basic vertices are: (a) absorption of a photon by an electron (b) emission of a photon by an electron (c) absorption of a photon by a positron (d) emission of a photon by a positron (e) electron-positron pair creation by a photon and (f) electron-positron pair annihilation into a photon. They are altogether six vertices as shown in Fig. 5.2. A single vertex is not a valid physical process since energy and momentum cannot be simultaneously conserved. So, a Feynman diagram for a valid physical process should consist of two or more such vertices and satisfy the following conditions:

1. Energy and momentum are conserved at each vertex.
2. The incoming and outgoing lines in a diagram are real particles that obey the relativistic energy-momentum relation $E^{2}=p^{2} c^{2}+m^{2} c^{4}$.
3. The internal lines (lines in the intermediate stages) are called the virtual particles that do not satisfy the energy-momentum relation but they form an integral part of the Feynman diagrams.

Any Feynman diagram consisting only of a single vertex cannot satisfy simultaneously the conditions 1 and 2 and hence cannot represent a physical process.

An electron is usually represented by a solid line with an arrow in the direction of increasing time, a positron by a solid line with an arrow opposite to the direction of increasing time and a photon by a wavy line.


Figure 5.2: Feynman diagrams representing basic vertices in QED: (a) photon absorption by electron, (b) photon emission by electron, (c) photon absorption by positron, (d) photon emission by positron, (e) electron-positron pair creation, (f) electron-positron pair annihilation.

For a given process, draw all possible diagrams and write down the matrix element for each diagram and add them up. In other words, one has to compute the transition amplitude corresponding to each diagram and find the total amplitude merely by summing them. The total transition probability is obtained by finding the absolute square of the total amplitude. If the spins of the initial and final particles are not observed,
then a sum over the final spin states and an average over the initial spin states have to be taken.

## Review Questions

5.1 Explain how you can study the time evolution of the state vector in Schrödinger equation by defining a propagation kernel $K_{0}(\mathbf{2}, \mathbf{1})$. By introducing a perturbing potential $V$, show how the propagation kernel $K(\mathbf{2}, \mathbf{1})$ with the perturbing potential can be expressed as an integral equation and obtain a perturbation series for $K(\mathbf{2}, \mathbf{1})$ in terms of $K_{0}(\mathbf{2}, \mathbf{1})$.
5.2 Obtain the propagation kernel for the wave function obeying the Dirac equation and explain how Feynman's modification gave a new insight that led to the development of Feynman's positron theory.
5.3 Obtain an expression for the transition amplitude for an electron to go from one positive energy state to another due to an interaction with electromagnetic field. Deduce the first order and second order matrix elements for such a transition both in space-time representation and momentumenergy representation.
5.4 Explain why the momentum-energy representation is preferable for the study of various processes in quantum electrodynamics.

## Problems

5.1 It is said that a single vertex which corresponds to a first order matrix element is not a valid physical process. But the Coulomb scattering of an electron by a nucleus is represented by a single vertex and calculated by using the first order matrix element. Explain the contradiction.
5.2 Draw the Feynman diagrams to describe the Compton scattering. Show that there are two possible Feynman diagrams which contribute to the process. Write down the matrix element corresponding to each diagram.
5.3 Show that there are two Feynman diagrams that contribute to the electronelectron scattering process in the lowest order. Obtain the matrix element for the two diagrams.

## Solutions to Problems

5.1 The Coulomb scattering of an electron by a nucleus is represented by the Feynman diagram as given in Fig. 6.1. Please note that the wavy line in figure denotes a virtual photon and not a real photon. A single vertex does not represent a physical process only if all the particles at a single vertex are real. So, there is no contradiction and Fig. 6.1 represents a valid physical process and so the cross section can be calculated using the first order matrix element (5.47).
5.2 The two Feynman diagrams for the Compton scattering are given in Fig. 6.3. It is a second order process and the matrix elements can be written down using Eq. (5.50). For details, refer to Chapter 6.
5.3 The problem of electron-electron scattering is treated in Sec. 6.3. along with Feynman diagrams. The reader is referred to that section for writing down the Feynman amplitudes for the process.

## Chapter 6

## Quantum Electrodynamics

Quantum Electrodynamics (QED) deals with processes involving interaction of charged particles with electromagnetic field. In this chapter, we shall deal with lowest order Feynman diagrams that describe the processes involving electrons and photons. If we go to higher orders, we are faced with problems of infinities and we shall briefly discuss how to tame them by renormalization techniques in the next chapter.

There are several processes, of which we shall choose the following for discussion: (1) Rutherford scattering (2) Compton scattering (3) Electronelectron scattering (4) Electron-positron scattering (5) Electron-positron pair annihilation (6) Bremsstrahlung (7) Electron-positron pair production and (8) Muon pair production in electron-positron collision. As Schwinger once remarked, Feynman has indeed taken the quantum elecrodynamical calculations "to the masses".

### 6.1 Rutherford Scattering

Consider the high energy scattering of an electron by the Coulomb potential of a nucleus of charge $Z e$ as depicted in Fig. 6.1. The potential scattering causes only a three-momentum transfer and no energy transfer. The matrix element for the process can be written down ${ }^{1}$ using Eq. (5.47).

$$
\begin{equation*}
\mathscr{M}_{1}=-i e(2 \pi)^{4} \bar{\psi}\left(\boldsymbol{p}_{f}\right) \phi(\mathbf{q}) \psi\left(\boldsymbol{p}_{i}\right) . \tag{6.1}
\end{equation*}
$$

[^39]

Figure 6.1: Coulomb scattering of electron by a nucleus of charge $Z e . \boldsymbol{p}_{i}$ and $\boldsymbol{p}_{f}$ denote the momenta of the incoming and scattered electron and $\boldsymbol{q}$, the threemomentum transfer to the electron. Since there is no energy transfer $\left|\boldsymbol{p}_{\boldsymbol{i}}\right|=$ $\left|\boldsymbol{p}_{f}\right|=p$ and $q^{2}=4 p^{2} \sin ^{2} \frac{\theta}{2}$.
where $\mathbf{q}=\mathbf{p}_{f}-\mathbf{p}_{i}$ is the four-momentum transfer $\left(q_{0}, \boldsymbol{q}\right)$ to the electron in the Coulomb interaction $V=Z e / r$ with a nucleus of charge $Z e, r$ being the distance measured from the centre of the nucleus. From Eq. (5.44), we obtain ${ }^{2}$

$$
\begin{align*}
\phi(\mathbf{q}) & =\frac{1}{(2 \pi)^{4}} \int e^{i \mathbf{q} \cdot \mathbf{x}_{\gamma_{\mu}} A_{\mu}(\mathbf{x}) d^{4} x} \\
& =\frac{1}{(2 \pi)^{4}} \int_{-\infty}^{\infty} e^{i q_{0} x_{0}} d x_{0} \int e^{-i \boldsymbol{q} \cdot \boldsymbol{x}} \gamma_{0} V(\boldsymbol{x}) d^{3} x \tag{6.2}
\end{align*}
$$

since $A_{0}(\mathbf{x})=V(\boldsymbol{x})=V(\boldsymbol{r})=\frac{Z e}{r}$ and $A_{k}(\mathbf{x})=0, k=1,2,3$ for a Coulomb interaction. The integrals in Eq. (6.2) can be evaluated using spherical coordinates and the results are given below.

$$
\begin{align*}
\int_{-\infty}^{\infty} e^{i q_{0} x_{0}} d x_{0} & =2 \pi \delta\left(q_{0}\right)  \tag{6.3}\\
\int e^{-i \boldsymbol{q} \cdot \boldsymbol{x}} V(\boldsymbol{x}) d^{3} x & =Z e \int e^{-i \boldsymbol{q} \cdot \boldsymbol{r}} \frac{1}{r} d^{3} r=\frac{4 \pi Z e}{q^{2}} \tag{6.4}
\end{align*}
$$

Substituting Eqs. (6.3) and (6.4) in Eq. (6.2), we get

$$
\begin{equation*}
\phi(\mathbf{q})=\frac{1}{(2 \pi)^{3}} \frac{4 \pi Z e}{q^{2}} \gamma_{0} \delta\left(q_{0}\right) . \tag{6.5}
\end{equation*}
$$

[^40]Substituting this $\phi(\mathbf{q})$ in Eq. (6.1) and taking its absolute square, we obtain the transition probability.

$$
\begin{equation*}
\left|\mathscr{M}_{1}\right|^{2}=(2 \pi)^{2}\left(\frac{4 \pi Z e^{2}}{q^{2}}\right)^{2}\left|\bar{\psi}\left(\boldsymbol{p}_{f}\right) \gamma_{0} \psi\left(\boldsymbol{p}_{i}\right)\right|^{2}\left[\delta\left(q_{0}\right)\right]^{2} \tag{6.6}
\end{equation*}
$$

where the square of the delta function requires some interpretation. Writing

$$
\begin{equation*}
\left[\delta\left(q_{0}\right)\right]^{2}=\delta\left(q_{0}\right) \delta(0) \tag{6.7}
\end{equation*}
$$

the first delta function can be interpreted to mean the conservation of energy and the second delta function $\delta(0)$ can be interpreted to denote the interaction time as shown below ${ }^{3}$.

$$
\begin{equation*}
\delta(0)=\lim _{T \rightarrow \infty} \lim _{x \rightarrow 0} \frac{1}{2 \pi} \int_{-T / 2}^{+T / 2} e^{i x t} d t=\lim _{T \rightarrow \infty} \frac{T}{2 \pi}, \tag{6.8}
\end{equation*}
$$

where $T$ is interpreted as the time during which interaction takes place. From Eqs. (6.6) - (6.8), we obtain the transition probability per unit time to a group of final states in the continuum with energy $E_{f}$.

$$
\begin{equation*}
w_{i \rightarrow f}=2 \pi\left(\frac{4 \pi Z e^{2}}{q^{2}}\right)^{2}\left|\bar{\psi}\left(\boldsymbol{p}_{f}\right) \gamma_{0} \psi\left(\boldsymbol{p}_{i}\right)\right|^{2} \delta\left(E_{i}-E_{f}\right) \rho_{f} \tag{6.9}
\end{equation*}
$$

where $\rho_{f}$ denotes the density of final states.
If we are considering an initial beam of unpolarized electron and not interested in the final spin states of the electron, then a sum over the final spin states and an average over the initial spin states have to be taken. For this, we use Eq. (2.75).

$$
\begin{equation*}
\bar{\sum}\left|\bar{\psi}\left(\boldsymbol{p}_{f}\right) \gamma_{0} \psi\left(\boldsymbol{p}_{i}\right)\right|^{2}=\frac{1}{2} \operatorname{Tr}\left(\gamma_{0}\left(p_{i}+m\right) \tilde{\gamma}_{0}\left(\not p_{f}+m\right)\right), \tag{6.10}
\end{equation*}
$$

where $\tilde{\gamma_{0}}=\gamma_{0} \gamma_{0}^{\dagger} \gamma_{0}=\gamma_{0}$. and $\not p=\gamma_{0} E-\gamma \cdot \boldsymbol{p}$. Since the trace ${ }^{4}$ of a

[^41]product of an odd number of $\gamma$ matrices vanishes, we get ${ }^{5}$
\[

$$
\begin{align*}
\overline{\sum \mid}\left|\bar{\psi}\left(\boldsymbol{p}_{f}\right) \gamma_{0} \psi\left(\boldsymbol{p}_{i}\right)\right|^{2} & =\frac{1}{2} \operatorname{Tr}(\underbrace{\gamma_{0} \not p_{i}} \gamma_{0} \not p_{f}+m^{2}) \\
& =\frac{1}{2} \operatorname{Tr}\left(\left(-\not p_{i} \gamma_{0}+2 E_{i}\right) \gamma_{0} \not p_{f}+m^{2}\right) \\
& =\frac{1}{2} \operatorname{Tr}\left(-\not p_{i} \not p_{f}+2 E_{i} \gamma_{0} \not p_{f}+m^{2}\right) \\
& =-2 \mathbf{p}_{i} \cdot \mathbf{p}_{f}+4 E_{i} E_{f}+2 m^{2} . \tag{6.11}
\end{align*}
$$
\]

Expanding the scalar product of the four-vectors $\mathbf{p}_{i} \cdot \mathbf{p}_{f}=E_{i} E_{f}-\boldsymbol{p}_{i} \cdot \boldsymbol{p}_{f}$ and rearranging, we get

$$
\begin{equation*}
\bar{\sum}\left|\bar{\psi}\left(\boldsymbol{p}_{f}\right) \gamma_{0} \psi\left(\boldsymbol{p}_{i}\right)\right|^{2}=2 E_{i} E_{f}+2 p_{i} p_{f} \cos \theta+2 m^{2}, \tag{6.12}
\end{equation*}
$$

where $\theta$ denotes the angle of scattering. In the potential scattering that we are considering, there is no energy transfer but there is only a momentum transfer. $E_{i}=E_{f}=E$ and $\left|\boldsymbol{p}_{i}\right|=\left|\boldsymbol{p}_{f}\right|=p$. Hence

$$
\begin{align*}
\bar{\sum}\left|\bar{\psi}\left(\boldsymbol{p}_{f}\right) \gamma_{0} \psi\left(\boldsymbol{p}_{i}\right)\right|^{2} & =2 E^{2}+2 p^{2} \cos \theta+2 m^{2} \\
& =4 E^{2}-2 p^{2}(1-\cos \theta), \text { since } m^{2}=E^{2}-p^{2} \\
& =4 E^{2}-4 p^{2} \sin ^{2}(\theta / 2) \\
& =4 E^{2}\left(1-v^{2} \sin ^{2}(\theta / 2)\right), \tag{6.13}
\end{align*}
$$

since ${ }^{6} v=p / E$ is the velocity of the incident electron.
The density of final states $\rho_{f}$ is given by

$$
\begin{equation*}
\rho_{f}=\frac{d n}{d E_{f}}=\frac{d^{3} p_{f}}{(2 \pi)^{3} d E_{f}}=\frac{p_{f}^{2} d p_{f} d \Omega}{(2 \pi)^{3} d E_{f}}=\frac{p_{f} E_{f} d \Omega}{(2 \pi)^{3}} . \tag{6.14}
\end{equation*}
$$

The last step is obtained using the relation $p_{f} d p_{f}=E_{f} d E_{f}$.
Substituting (6.13) and (6.14) in Eq. (6.9) and introducing the normalization factor $\sqrt{2 E}$ for each fermion wave function such that $N=2 E$ and remembering that $E_{f}=E$ and $p_{f}=p$, we get the transition probability per unit time (the transition rate).

$$
\begin{equation*}
w_{i \rightarrow f}=\frac{2 \pi}{N^{2}}\left(\frac{4 \pi Z e^{2}}{q^{2}}\right)^{2} 4 E^{2}\left(1-v^{2} \sin ^{2} \frac{\theta}{2}\right) \frac{p E}{(2 \pi)^{3}} . \tag{6.15}
\end{equation*}
$$

[^42]Since the incident wave is normalized to 1 particle per unit volume, the incident flux is $v$.

$$
\text { Cross Section }=\frac{\text { Transition Rate }}{\text { Incident Flux }}
$$

Thus

$$
\begin{equation*}
\frac{d \sigma}{d \Omega}=\frac{2 \pi}{v}\left(\frac{4 \pi Z e^{2}}{q^{2}}\right)^{2}\left(1-v^{2} \sin ^{2} \frac{\theta}{2}\right) \frac{p E}{(2 \pi)^{3}} \tag{6.16}
\end{equation*}
$$

Substituting $q^{2}=4 p^{2} \sin ^{2}(\theta / 2)$ in Eq. (6.16) and simplifying, we get

$$
\begin{equation*}
\frac{d \sigma}{d \Omega}=\frac{1}{v^{2}} \frac{Z^{2} e^{4}}{4 p^{2} \sin ^{4}(\theta / 2)}\left(1-v^{2} \sin ^{2} \frac{\theta}{2}\right) . \tag{6.17}
\end{equation*}
$$

Equation (6.17) is the Mott scattering cross section. In the non-relativistic limit, since $v=\frac{p}{E} \ll 1$, the cross section (6.17) reduces to the Rutherford scattering cross section.

$$
\begin{equation*}
\frac{d \sigma}{d \Omega}=\frac{1}{v^{2}} \frac{Z^{2} e^{4}}{4 p^{2} \sin ^{4}(\theta / 2)} \tag{6.18}
\end{equation*}
$$

The relativistic effect is to introduce the additional factor $\left(1-v^{2} \sin ^{2}(\theta / 2)\right)$ to the Rutherford scattering formula.

### 6.1.1 Fermi's golden rule

Above, we have deduced the cross section (6.17) from first principles. But it is found convenient to use Fermi's golden rule to obtain the transition probability per second (transition rate) and differential cross section.

$$
\begin{align*}
\text { Transition rate } & =2 \pi(\Pi N)^{-1}|\overline{\mathcal{M}}|^{2} \rho_{f}  \tag{6.19}\\
d \sigma=\frac{\text { Transition rate }}{\text { Incident flux }} & =\frac{2 \pi}{v}(\Pi N)^{-1}|\overline{\mathcal{M}}|^{2} \rho_{f} \tag{6.20}
\end{align*}
$$

where $(\Pi N)^{-1}$ denotes Feynman's normalization factors, $|\overline{\mathcal{M}}|^{2}$ stands for the square of the transition matrix element, obtained after summing over the final spin states and averaging over the initial spin states and $\rho_{f}$ is
the density of final states ${ }^{7}$. In the present case,

$$
\begin{align*}
(\Pi N)^{-1} & =\frac{1}{4 E^{2}},  \tag{6.21}\\
|\overline{\mathcal{M}}|^{2} & \left.=\frac{1}{2} \sum_{\text {spins }}\left|\left\langle u_{f}\right| \mathcal{O}\right| u_{i}\right\rangle\left.\right|^{2}, \quad \text { with } \mathcal{O}=\frac{4 \pi z e^{2}}{q^{2}} \gamma_{0},  \tag{6.22}\\
\rho_{f} & =\frac{p E}{(2 \pi)^{3}} d \Omega . \tag{6.23}
\end{align*}
$$

Substituting (6.21) - (6.23) into Eq. (6.20), we obtain the Mott differential cross section (6.17).

In the consideration of the other quantum electrodynamic processes, discussed below in this chapter, we shall use Fermi's golden rule for calculating the differential cross sections.

### 6.2 Compton Scattering

The scattering of a photon by an electron is known as Compton scattering, for which the relativistic cross section was first deduced by Klein and Nishina ${ }^{8}$. The potential corresponding to the absorption of a photon is given by

$$
\begin{equation*}
A_{\mu}=\left(\frac{4 \pi e^{2}}{2 \omega}\right)^{1 / 2} \epsilon_{\mu} e^{i \mathbf{k} \cdot \mathbf{x}} \tag{6.24}
\end{equation*}
$$

For the emission of photon, one has to take $A_{\mu}^{*}$, the complex conjugate of (6.24). The potential $A_{\mu}$ in (6.24) is normalized to 1 photon per unit volume which is not relativistically invariant. For relativistic invariance, the potential should be normalized to $2 \omega$ photons per unit volume. So, the relativistically invariant potential for the absorption of a photon can be written as

$$
\begin{equation*}
A_{\mu}=\left(4 \pi e^{2}\right)^{1 / 2} \epsilon_{\mu} e^{i \mathbf{k} \cdot \mathbf{x}} \tag{6.25}
\end{equation*}
$$

To obtain the correct transition probability in a given coordinate system, it is necessary to reinsert a factor $(2 \omega)^{-1}$ for each photon in the initial and final states. This becomes a part of the normalization factor $\Pi N$ which contains a similar factor for each electron in the initial and final states.

[^43]In momentum representation, the amplitude to absorb (emit) a photon of polarization $\epsilon_{\mu}$ is

$$
-i\left(4 \pi e^{2}\right)^{1 / 2} k
$$

The polarization vector $\boldsymbol{\epsilon}$ is perpendicular to the propagation vector $\boldsymbol{k}$.

$$
\boldsymbol{\epsilon} \cdot \boldsymbol{k}=0 .
$$

### 6.2.1 Kinematics

As a first approximation, consider the electrons to be free and choose the laboratory system, in which the initial electron is at rest. A photon of four-momentum $\mathbf{k}_{1}$ is incident along the x -axis on an electron at rest and the scattered photon makes an angle $\theta$ with the incident direction and carries a four-momentum $\mathbf{k}_{2}$ as shown in Fig. 6.2. The electron recoils at an angle $\phi$.


Figure 6.2: Compton scattering in the rest frame of the initial electron.
We shall represent the incoming and outgoing photons by the potentials

$$
\begin{equation*}
A_{1 \mu}=\epsilon_{1 \mu} e^{-i \mathbf{k}_{1} \cdot \mathbf{x}}, \quad A_{2 \mu}=\epsilon_{2 \mu} e^{-i \mathbf{k}_{2} \cdot \mathbf{x}} \tag{6.26}
\end{equation*}
$$

and the initial and final electron states by the wave functions

$$
\begin{equation*}
\psi_{1}=u_{1} e^{-i \mathbf{p}_{1} \cdot \mathbf{x}}, \quad \psi_{2}=u_{2} e^{-i \mathbf{p}_{2} \cdot \mathbf{x}} \tag{6.27}
\end{equation*}
$$

The electron spinors obey the Dirac equation

$$
\begin{equation*}
\not p_{1} u_{1}=m u_{1}, \quad \not p_{2} u_{2}=m u_{2}, \tag{6.28}
\end{equation*}
$$

such that

$$
\begin{equation*}
\mathbf{p}_{1} \cdot \mathbf{p}_{1}=m^{2}, \quad \mathbf{p}_{2} \cdot \mathbf{p}_{2}=m^{2} \tag{6.29}
\end{equation*}
$$

For a photon, the energy and momentum are both equal to frequency in natural units $(\hbar=c=1)$ such that

$$
\begin{equation*}
\mathbf{k}_{1} \cdot \mathbf{k}_{1}=0, \quad \mathbf{k}_{2} \cdot \mathbf{k}_{2}=0 \tag{6.30}
\end{equation*}
$$

The polarization four-vector has no time component but its space components are normal to its momentum vector. Hence

$$
\begin{equation*}
\varepsilon_{1} \cdot \mathbf{k}_{1}=0, \quad \varepsilon_{2} \cdot \mathbf{k}_{2}=0 \tag{6.31}
\end{equation*}
$$

Conservation of energy and momentum is given by

$$
\begin{equation*}
\mathbf{p}_{1}+\mathbf{k}_{1}=\mathbf{p}_{2}+\mathbf{k}_{2} \quad \text { or } \quad \not p_{1}+\not x_{1}=\not p_{2}+\not \not{ }_{2} \tag{6.32}
\end{equation*}
$$

One of the angles, either $\theta$ or $\phi$, is sufficient to determine the remaining quantities. Choosing the rest system of the initial electron, we get

$$
\begin{align*}
\not p_{2} & =\not p_{1}+\not k_{1}-\not \phi_{2} ; \\
\not p_{2}^{2} & =\left(\not p_{1}+\not \phi_{1}-\not \phi_{2}\right)\left(\not p_{1}+\not b_{1}-\not p_{2}\right) \\
& =\mathbf{p}_{1}^{2}+\mathbf{k}_{1}^{2}+\mathbf{k}_{2}^{2}+2 \mathbf{p}_{1} \cdot \mathbf{k}_{1}-2 \mathbf{p}_{1} \cdot \mathbf{k}_{2}-2 \mathbf{k}_{1} \cdot \mathbf{k}_{2} \\
& =m^{2}+2 m \omega_{1}-2 m \omega_{2}-2 \omega_{1} \omega_{2}(1-\cos \theta) \tag{6.33}
\end{align*}
$$

Since $\not p_{2}^{2}$ is also equal to

$$
\begin{equation*}
\not p_{2}^{2}=\mathbf{p}_{2} \cdot \mathbf{p}_{2}=m^{2} \tag{6.34}
\end{equation*}
$$

we obtain the following relation from Eqs. (6.33) and (6.34).

$$
\begin{align*}
m\left(\omega_{1}-\omega_{2}\right) & =\omega_{1} \omega_{2}(1-\cos \theta) \\
\frac{m}{\omega_{2}}-\frac{m}{\omega_{1}} & =1-\cos \theta \tag{6.35}
\end{align*}
$$

This is the well-known formula for the Compton shift in frequency or wavelength.

### 6.2.2 Transition matrix element

Compton scattering is a process in which an electron of four-momentum $\mathbf{p}_{1}$ and photon of four-momentum $\mathbf{k}_{1}$ in the initial state make a transition to a final state of electron of four-momentum $\mathbf{p}_{2}$ and photon of fourmomentum $\mathbf{k}_{2}$. In the lowest order, it can be envisaged as a second-order process in which an electron first absorbs a photon of four-momentum


Figure 6.3: Lowest order Feynman diagrams for Compton scattering.
$\mathbf{k}_{1}$ and then emits a photon of four-momentum $\mathbf{k}_{2}$. This process is indistinguishable from a process in which the electron first emits a photon of four-momentum $\mathbf{k}_{2}$ and then absorbs a photon of four-momentum $\mathbf{k}_{1}$. These two modes can be represented by Feynman diagrams (a) and (b) as shown in Fig. 6.3.

Using the Feynman rules, the matrix elements for the diagrams (a) and (b) can be written as

$$
\begin{align*}
& \mathcal{M}_{\mathrm{a}}=-i 4 \pi e^{2}\left(\bar{u}_{2} \not_{2} \frac{1}{\not p_{1}+\not \phi_{1}-m} \phi_{1} u_{1}\right),  \tag{6.36}\\
& \mathcal{M}_{\mathrm{b}}=-i 4 \pi e^{2}\left(\bar{u}_{2} \phi_{1} \frac{1}{\not p_{1}-\not \phi_{2}-m} \xi_{2} u_{1}\right) . \tag{6.37}
\end{align*}
$$

Taking the sum of these two contributions, we get

$$
\begin{equation*}
\mathcal{M}=\mathcal{M}_{\mathrm{a}}+\mathcal{M}_{\mathrm{b}}=-i 4 \pi e^{2}\left\{\bar{u}_{2}\left(\mathcal{O}_{\mathrm{a}}+\mathcal{O}_{\mathrm{b}}\right) u_{1}\right\}, \tag{6.38}
\end{equation*}
$$

with

$$
\begin{align*}
\mathcal{O}_{\mathrm{a}} & =\xi_{2} \frac{1}{\not p_{1}+\not k_{1}-m} \xi_{1}=\frac{\not \phi_{2}\left(\not p_{1}+\not k_{1}+m\right) \xi_{1}}{\left(\not p_{1}+\not 1_{1}\right)^{2}-m^{2}},  \tag{6.39}\\
\mathcal{O}_{\mathrm{b}} & =\not \phi_{1} \frac{1}{\not p_{1}-\not k_{2}-m} \phi_{2}=\frac{\not \phi_{1}\left(\not p_{1}-\not \phi_{2}+m\right) \xi_{2}}{\left(\not p_{1}-\not \phi_{2}\right)^{2}-m^{2}} . \tag{6.40}
\end{align*}
$$

Here, we make an interesting observation. The matrix element $\mathcal{M}_{\mathrm{a}}$ goes over into $\mathcal{M}_{\mathrm{b}}$ by making the substitution

$$
\begin{equation*}
\mathbf{k}_{1} \longrightarrow-\mathbf{k}_{2} \quad \varepsilon_{1} \longrightarrow \varepsilon_{2} . \tag{6.41}
\end{equation*}
$$

and the total matrix element (the total transition amplitude) is invariant by the substitution (6.41). Such symmetries are known as crossing symmetries.

We can simplify the expressions (6.39) and (6.40) by choosing a coordinate system in which the initial electron is at rest. The denominators now simplify to

$$
\begin{align*}
& \left(\not p_{1}+\not \not{ }_{1}\right)^{2}-m^{2}=\not p_{1}^{2}+\not \not{ }_{1}^{2}+2 \not p_{1} \not{ }_{1}-m^{2} \\
& =\mathbf{p}_{1}^{2}+\mathbf{k}_{1}^{2}+2 \mathbf{p}_{1} \cdot \mathbf{k}_{1}-m^{2}=2 m \omega_{1},  \tag{6.42}\\
& \left(\not p_{1}-\not \nmid_{2}\right)^{2}-m^{2}=\not p_{1}^{2}+\not \nmid_{2}^{2}-2 \not p_{1} \not \phi_{2}-m^{2} \\
& =\mathbf{p}_{1}^{2}+\mathbf{k}_{2}^{2}-2 \mathbf{p}_{1} \cdot \mathbf{k}_{2}-m^{2}=-2 m \omega_{2}, \tag{6.43}
\end{align*}
$$

since

$$
\mathbf{p}_{1}^{2}=m^{2}, \quad \mathbf{k}_{1}^{2}=0, \quad 2 \mathbf{p}_{1} \cdot \mathbf{k}_{1}=2 m \omega_{1}, \quad \mathbf{k}_{2}^{2}=0, \quad 2 \mathbf{p}_{1} \cdot \mathbf{k}_{2}=2 m \omega_{2}
$$

The numerators also can be similarly simplified.

$$
\begin{align*}
& \phi_{2} \not b_{1} \phi_{1}=-\phi_{2} \phi_{1} \not p_{1}+2 \mathbf{p}_{1} \cdot \varepsilon_{1}=-\phi_{2} \phi_{1} \not p_{1},  \tag{6.44}\\
& \phi_{1} b_{1} \phi_{2}=-\phi_{1} \phi_{2} \not p_{1}+2 \mathbf{p}_{1} \cdot \varepsilon_{2}=-\phi_{1} \phi_{2} \not p_{1}, \tag{6.45}
\end{align*}
$$

since $\mathbf{p}_{1}$ has only a time component whereas $\varepsilon_{1}$ and $\varepsilon_{2}$ have only space components such that $\mathbf{p}_{1} \cdot \varepsilon_{1}=0$ and $\mathbf{p}_{1} \cdot \varepsilon_{2}=0$. Further when taken between the electron spinors, they yield

$$
\begin{align*}
& \bar{u}_{2} \phi_{2} p_{1} \phi_{1} u_{1}=-\bar{u}_{2} \phi_{2} \phi_{1} \phi_{1} u_{1}=-m \bar{u}_{2} \phi_{2} \phi_{1} u_{1},  \tag{6.46}\\
& \bar{u}_{2} \phi_{1} b_{1} \phi_{2} u_{1}=-\bar{u}_{2} \phi_{1} \phi_{2} \phi_{1} u_{1}=-m \bar{u}_{2} \phi_{1} \phi_{2} u_{1}, \tag{6.47}
\end{align*}
$$

since $\not p_{1} u_{1}=m u_{1}$. Substituting (6.42) - (6.47) into Eqs. (6.39) and (6.40) and simplifying, we get

$$
\begin{align*}
& \tilde{u}_{2} \mathcal{O}_{\mathrm{a}} u_{1}=\frac{\bar{u}_{2} \not_{2} \not k_{1} \xi_{1} u_{1}}{2 m \omega_{1}}=-\frac{\bar{u}_{2} \xi_{2} \not_{1} \not k_{1} u_{1}}{2 m \omega_{1}} ;  \tag{6.48}\\
& \bar{u}_{2} \mathcal{O}_{\mathrm{b}} u_{1}=\frac{\bar{u}_{2} \xi_{1} k_{2} \xi_{2} u_{1}}{2 m \omega_{2}}=-\frac{\bar{u}_{2} \xi_{1} \xi_{2} k_{2} u_{1}}{2 m \omega_{2}} . \tag{6.49}
\end{align*}
$$

The last step in Eqs. (6.48) is obtained by using the commutation relation $\not \phi_{1} \xi_{1}=-\xi_{1} \not k_{1}+2 \mathbf{k}_{1} \cdot \varepsilon_{1}$. The scalar product of four vectors vanishes since $\varepsilon_{1}$ has no time component and the polarization three vector is normal to the photon three-momentum. A similar argument leads also to the last step in Eq. (6.49).

### 6.2.3 Transition probability

The absolute square of the transition matrix element $\mathcal{M}$ yields the transition probability. If the spin states of the electron in the initial and final states are not observed, then a sum over the final spin states and an average over the initial spin states have to be done. This is done by using the trace techniques discussed earlier.

$$
\begin{equation*}
|\overline{\mathcal{M}}|^{2}=16 \pi^{2} e^{4} \bar{\sum}\left|\bar{u}_{2} \mathcal{O} u_{1}\right|^{2}, \tag{6.50}
\end{equation*}
$$

where ${ }^{9}$

$$
\begin{align*}
\mathcal{O} & =\mathcal{O}_{\mathrm{a}}+\mathcal{O}_{\mathrm{b}} \\
& =\frac{\phi_{2} \xi_{1} \not \phi_{1}}{2 m \omega_{1}}+\frac{\xi_{1} \xi_{2} \not \phi_{2}}{2 m \omega_{2}} . \tag{6.51}
\end{align*}
$$

Using the trace techniques, we get

$$
\begin{align*}
\bar{\sum}\left|\bar{u}_{2} \mathcal{O} u_{1}\right|^{2} & =\frac{1}{2} \operatorname{Tr}\left[\mathcal{O}\left(\not p_{1}+m\right) \tilde{\mathcal{O}}\left(\not p_{2}+m\right)\right] \\
& =\frac{1}{2} \operatorname{Tr}\left[\left(\mathcal{O}_{\mathrm{a}}+\mathcal{O}_{\mathrm{b}}\right)\left(\not p_{1}+m\right)\left(\tilde{\mathcal{O}}_{\mathrm{a}}+\tilde{\mathcal{O}}_{\mathrm{b}}\right)\left(\not p_{2}+m\right)\right] \\
& =\frac{1}{2}\left[T_{\mathrm{aa}}+T_{\mathrm{ab}}+T_{\mathrm{ba}}+T_{\mathrm{bb}}\right] \tag{6.52}
\end{align*}
$$

where

$$
\begin{equation*}
\tilde{\mathcal{O}}=\gamma_{0} \mathcal{O}^{\dagger} \gamma_{0}, \quad \tilde{\mathcal{O}}_{\mathrm{a}}=\gamma_{0} \mathcal{O}_{\mathrm{a}}^{\dagger} \gamma_{0}, \quad \tilde{\mathcal{O}}_{\mathrm{b}}=\gamma_{0} \mathcal{O}_{\mathrm{b}}^{\dagger} \gamma_{0} \tag{6.53}
\end{equation*}
$$

and

$$
\begin{align*}
& T_{\mathrm{aa}}=\operatorname{Tr}\left[\mathcal{O}_{\mathrm{a}}\left(\not p_{1}+m\right) \tilde{\mathcal{O}}_{\mathrm{a}}\left(\not p_{2}+m\right)\right] ;  \tag{6.54}\\
& T_{\mathrm{ab}}=\operatorname{Tr}\left[\mathcal{O}_{\mathrm{a}}\left(\not p_{1}+m\right) \tilde{\mathcal{O}}_{\mathrm{b}}\left(\not p_{2}+m\right)\right] ;  \tag{6.55}\\
& T_{\mathrm{ba}}=\operatorname{Tr}\left[\mathcal{O}_{\mathrm{b}}\left(\not p_{1}+m\right) \tilde{\mathcal{O}}_{\mathrm{a}}\left(\not p_{2}+m\right)\right] ;  \tag{6.56}\\
& T_{\mathrm{bb}}=\operatorname{Tr}\left[\mathcal{O}_{\mathrm{b}}\left(\not p_{1}+m\right) \tilde{\mathcal{O}}_{\mathrm{b}}\left(\not p_{2}+m\right)\right] . \tag{6.57}
\end{align*}
$$

[^44]Writing explicitly the operators $\mathcal{O}_{\mathrm{a}}, \mathcal{O}_{\mathrm{b}}, \tilde{\mathcal{O}}_{\mathrm{a}}$ and $\tilde{\mathcal{O}}_{\mathrm{b}}$,

$$
\begin{align*}
& \mathcal{O}_{\mathrm{a}}=\frac{1}{2 m \omega_{1}} \xi_{2} \xi_{1} k_{1},  \tag{6.58}\\
& \mathcal{O}_{\mathrm{b}}=\frac{1}{2 m \omega_{2}} \xi_{1} \xi_{2} \not k_{2} \text {, }  \tag{6.59}\\
& \tilde{\mathcal{O}}_{\mathrm{a}}=\frac{1}{2 m \omega_{1}} \gamma_{0}\left(\phi_{2} \phi_{1} \not \phi_{1}\right)^{\dagger} \gamma_{0}=\frac{1}{2 m \omega_{1}}\left(\not k_{1} \phi_{1} \phi_{2}\right) \text {, }  \tag{6.60}\\
& \tilde{\mathcal{O}}_{\mathrm{b}}=\frac{1}{2 m \omega_{2}} \gamma_{0}\left(\xi_{1} \xi_{2} \not \phi_{2}\right)^{\dagger} \gamma_{0}=\frac{1}{2 m \omega_{2}}\left(\not \xi_{2} \phi_{2} \phi_{1}\right) \text {, } \tag{6.61}
\end{align*}
$$

we can evaluate the traces $T_{\mathrm{aa}}, T_{\mathrm{ab}}, T_{\mathrm{ba}}$ and $T_{\mathrm{bb}}$. These have been given as Problem 4.2 and its solution given at the end of this chapter. Here, we only give the final results.

$$
\begin{align*}
& T_{\mathrm{aa}}=\frac{2}{m \omega_{1}}\left\{m \omega_{2}+2\left(\mathbf{k}_{1} \cdot \boldsymbol{\varepsilon}_{2}\right)^{2}\right\}  \tag{6.62}\\
& T_{\mathrm{bb}}=\frac{2}{m \omega_{2}}\left\{m \omega_{1}-2\left(\mathbf{k}_{2} \cdot \boldsymbol{\varepsilon}_{1}\right)^{2}\right\}  \tag{6.63}\\
& T_{\mathrm{ab}}=2\left\{2\left(\varepsilon_{1} \cdot \boldsymbol{\varepsilon}_{2}\right)^{2}-1\right\}-\frac{2}{m \omega_{1}}\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)^{2}+\frac{2}{m \omega_{2}}\left(\mathbf{k}_{2} \cdot \boldsymbol{\varepsilon}_{1}\right)^{2} .  \tag{6.64}\\
& T_{\mathrm{ba}}=T_{\mathrm{ab}} . \tag{6.65}
\end{align*}
$$

The term $T_{\mathrm{bb}}$ can be obtained from $T_{\mathrm{aa}}$ by using the crossing symmetry $\left(\mathbf{k}_{1} \longrightarrow-\mathbf{k}_{2}, \varepsilon_{1} \longrightarrow \varepsilon_{2}, \omega_{1} \longrightarrow \omega_{2}\right.$ but the term $T_{\mathrm{ab}}$ is invariant under the transformation representing crossing symmetry and it is equal to $T_{\mathrm{ba}}$.

Substituting (6.62) - (6.65) into Eq. (6.52), we obtain

$$
\begin{equation*}
\bar{\sum}\left|\bar{u}_{2} \mathcal{O} u_{1}\right|^{2}=\frac{\omega_{2}}{\omega_{1}}+\frac{\omega_{1}}{\omega_{2}}+4\left(\varepsilon_{1} \cdot \varepsilon_{2}\right)^{2}-2 \tag{6.66}
\end{equation*}
$$

Consequently the matrix element square, given by Eq. (6.50) becomes

$$
\begin{equation*}
|\overline{\mathcal{M}}|^{2}=16 \pi^{2} e^{4}\left\{\frac{\omega_{2}}{\omega_{1}}+\frac{\omega_{1}}{\omega_{2}}+4\left(\varepsilon_{1} \cdot \varepsilon_{2}\right)^{2}-2\right\} . \tag{6.67}
\end{equation*}
$$

Transition probability per unit time $=2 \pi(\Pi N)^{-1}|\overline{\mathcal{M}}|^{2} \rho_{f}$,
where $\Pi N$ denotes the normalization factors that should be included for the incoming and outgoing particles and $\rho_{f}$, the density of final states. On dividing the transition probability per unit time (transition rate) by
incident flux, we obtain the cross section. Since the incident flux is $c$ which is unity in natural units, we get

$$
\begin{align*}
d \sigma & =\frac{2 \pi}{2 E_{1} 2 E_{2} 2 \omega_{1} 2 \omega_{2}}|\overline{\mathcal{M}}|^{2} \rho_{f} \\
& =\frac{2 \pi}{2 m 2 E_{2} 2 \omega_{1} 2 \omega_{2}}|\overline{\mathcal{M}}|^{2} \rho_{f} \tag{6.69}
\end{align*}
$$

since $E_{1}=m$ in the laboratory system in which the initial electron is at rest.

### 6.2.4 Density of final states

For the two-particle final states, the density of states is given by (vide Appendix B)

$$
\begin{equation*}
\rho_{f}=(2 \pi)^{-3} E_{a} E_{b} \frac{p_{a}^{3} d \Omega_{a}}{E p_{a}^{2}-E_{a}\left(\boldsymbol{p} \cdot \boldsymbol{p}_{a}\right)}, \tag{6.70}
\end{equation*}
$$

where $E=E_{a}+E_{b}$ and $\boldsymbol{p}=\boldsymbol{p}_{a}+\boldsymbol{p}_{b}$. Identifying the particle a with the outgoing photon and particle b with the scattered electron, we get

$$
\begin{array}{ll}
E_{a}=\omega_{2} ; \quad E_{b}=E_{2} ; \quad \begin{array}{l}
E=\omega_{2}+E_{2}=m+\omega_{1} . \\
\\
p_{a}=\omega_{2} ; \quad p_{b}=p_{2} ; \\
p=\omega_{2}+p_{2}=\omega_{1} .
\end{array}
\end{array}
$$

(since momentum is conserved)
Substituting these values in Eq. (6.70), we get

$$
\begin{align*}
\rho_{f} & =(2 \pi)^{-3} \omega_{2} E_{2} \frac{\omega_{2}^{3} d \Omega_{\omega}}{\left(\omega_{2}+E_{2}\right) \omega_{2}^{2}-\omega_{2}\left(\omega_{1} \omega_{2} \cos \theta\right)} \\
& =(2 \pi)^{-3} \frac{E_{2} \omega_{2}^{3} d \Omega_{\omega}}{\left(m+\omega_{1}\right) \omega_{2}-\omega_{1} \omega_{2} \cos \theta} \\
& =(2 \pi)^{-3} \frac{E_{2} \omega_{2}^{3} d \Omega_{\omega}}{m \omega_{1}} . \tag{6.71}
\end{align*}
$$

The last step is obtained using the relation (6.35) for the Compton frequency shift.

### 6.2.5 Klein-Nishina formula

Substituting Eqs. (6.67) and (6.71) into Eq. (6.69), we obtain the KleinNishina formula for the Compton scattering cross section.

$$
\begin{equation*}
\frac{d \sigma}{d \Omega_{\omega}}=\frac{e^{4}}{4 m^{2}}\left(\frac{\omega_{2}}{\omega_{1}}\right)^{2}\left[\frac{\omega_{2}}{\omega_{1}}+\frac{\omega_{1}}{\omega_{2}}+4\left(\varepsilon_{1} \cdot \varepsilon_{2}\right)^{2}-2\right] . \tag{6.72}
\end{equation*}
$$

For simplicity, we have derived the formula for real polarization vectors which corresponds to linear polarization. If the polarization vectors are complex, then one has to replace $\left(\varepsilon_{1} \cdot \varepsilon_{2}\right)^{2}$ by $\left(\varepsilon_{1} \cdot \varepsilon_{2}^{*}\right)^{2}$ in Eq. (6.72). If the polarizations of the incident and final photons are not observed, then the appropriate cross section is obtained by averaging the initial polarization states and summing over the final polarization states. The result of this operation is to replace $\left(\varepsilon_{1} \cdot \varepsilon_{2}\right)^{2}$ by $\frac{1}{2}\left(1+\cos ^{2} \theta\right)$, while the terms independent of $\varepsilon_{1}$ and $\varepsilon_{2}$ are multiplied by 2 .

$$
\begin{equation*}
\frac{d \sigma}{d \Omega_{\omega}}=\frac{e^{4}}{2 m^{2}}\left(\frac{\omega_{2}}{\omega_{1}}\right)^{2}\left[\frac{\omega_{2}}{\omega_{1}}+\frac{\omega_{1}}{\omega_{2}}-\sin ^{2} \theta\right] . \tag{6.73}
\end{equation*}
$$

## Total cross section

Equation (6.73) gives the differential cross section. To obtain the total cross section, we need to perform the angular integration. A word of caution is necessary. The quantity $\omega_{2}$ occurring in expression (6.73) depends on the scattering angle $\theta$ as given by Eq. (6.35). Multiplying Eq. (6.35) by $\omega_{1} / m$, we obtain

$$
\begin{align*}
& \frac{\omega_{1}}{\omega_{2}}=1+\frac{\omega_{1}}{m}(1-\cos \theta)=1+r(1-\cos \theta)  \tag{6.74}\\
& \frac{\omega_{2}}{\omega_{1}}=\frac{1}{1+\frac{\omega_{1}}{m}(1-\cos \theta)}=\frac{1}{1+r(1-\cos \theta)} \tag{6.75}
\end{align*}
$$

where $r=\frac{\omega_{1}}{m}$. The angular integration is performed taking into account the angular dependence of $\omega_{2}$ (vide solved problem 4.3) to obtain the total cross section.

$$
\begin{equation*}
\sigma_{\text {total }}=\frac{2 \pi e^{4}}{m^{2}}\left[\frac{1+r}{(1+2 r)^{2}}+\frac{2}{r^{2}}-\frac{2(1+r)-r^{2}}{2 r^{3}} \ln (1+2 r)\right] . \tag{6.76}
\end{equation*}
$$

## Non-relativistic limit

In the non-relativistic limit of photon-electron scattering, when the energy of the incident photon is very small, $\omega_{1} \ll m$, then $\omega_{2} \approx \omega_{1}$ and the scattering can be treated as elastic. In this limiting case, Eq. (6.73) simplifies to

$$
\begin{equation*}
\frac{d \sigma}{d \Omega_{\omega}}=\frac{e^{4}}{2 m^{2}}\left(2-\sin ^{2} \theta\right)=\frac{e^{4}}{2 m^{2}}\left(1+\cos ^{2} \theta\right) \tag{6.77}
\end{equation*}
$$

Performing the angular integration, we obtain the total cross section.

$$
\begin{equation*}
\sigma=\frac{e^{4}}{2 m^{2}}\left(4 \pi+\frac{4 \pi}{3}\right)=\frac{8 \pi e^{4}}{3 m^{2}}=0.665 \times 10^{-24} \mathrm{~cm}^{2} . \tag{6.78}
\end{equation*}
$$

This is the Thomson scattering formula and can be considered as the non-relativistic limit of Compton scattering for unpolarized photons.

## Extreme relativistic limit

In the other extreme relativistic limit when the incident photon energy is extremely large ( $\omega_{1} \gg m$ ), then Eq. (6.35) reduces to

$$
\omega_{2} \approx \frac{m}{1-\cos \theta}
$$

and the differential cross section (6.73) becomes

$$
\begin{equation*}
\frac{d \sigma}{d \Omega_{\omega}}=\frac{e^{4}}{2 m^{2}} \frac{\omega_{2}}{\omega_{1}}=\frac{e^{4}}{2 m \omega_{1}(1-\cos \theta)}=\frac{e^{4}}{4 m \omega_{1} \sin ^{2}(\theta / 2)} \tag{6.79}
\end{equation*}
$$

The total cross section in the extreme relativistic limit follows from Eq. (6.76).

$$
\begin{equation*}
\sigma=\frac{\pi e^{4}}{m \omega_{1}}\left\{\ln \frac{2 \omega_{1}}{m}+\frac{1}{2}\right\} \tag{6.80}
\end{equation*}
$$

### 6.3 Electron-electron scattering

Electron-electron scattering in the lowest order can be described by an exchange of one photon. Since electrons are identical particles, it is not possible to distinguish the two electrons in the final state. So, the scattering process is represented by two Feynman diagrams (a) and (b) as shown in Fig. 6.4. In Fig. 6.4(b), the four-momenta of the outgoing electrons are exchanged.

The matrix element for the process can easily be written down in the momentum representation by using the Feynman rules. For the sake of brevity, we denote the spinors $u\left(\mathbf{p}_{1}\right), u\left(\mathbf{p}_{2}\right) \cdots$ by $u_{1}, u_{2}, \cdots$.

$$
\begin{equation*}
\mathcal{M}=i 4 \pi e^{2}\left[\frac{\left(\bar{u}_{4} \gamma_{\mu} u_{2}\right)\left(\bar{u}_{3} \gamma_{\mu} u_{1}\right)}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{2}}-\frac{\left(\bar{u}_{3} \gamma_{\mu} u_{2}\right)\left(\bar{u}_{4} \gamma_{\mu} u_{1}\right)}{\left(\mathbf{p}_{1}-\mathbf{p}_{4}\right)^{2}}\right], \tag{6.81}
\end{equation*}
$$

where summation over $\mu$ is implied. Thus, the longitudinal and transverse waves of the photon are taken into account in a relativistically invariant


Figure 6.4: Electron-electron scattering. $\mathbf{p}_{1}$ and $\mathbf{p}_{2}$ denote the four-momenta of the electrons before collision and $\mathbf{p}_{3}$ and $\mathbf{p}_{4}$ after collision. Due to the indistinguishability of the two electrons in the final state, the exchange diagram (b) should also be considered.
way. The first term in Eq. (6.81) corresponds to the diagram (a) and the second term corresponds to the exchange diagram (b). The negative sign between the two terms arises from the antisymmetry of the two-electron wave function. To find the transition probability, we have to find the absolute square ${ }^{10}$ of the matrix element (6.81). If the electron spins are not observed, then a sum over the final spin states and an average over the initial spin states have to be taken. This introduces a factor $1 / 4$.

$$
\begin{align*}
|\overline{\mathcal{M}}|^{2}= & \frac{1}{4} \mathcal{M}^{\prime} \mathcal{M}^{\dagger} \\
= & 4 \pi^{2} e^{4} \sum_{\text {spins }}\left[\frac{\left(\bar{u}_{4} \gamma_{\mu} u_{2}\right)\left(\bar{u}_{3} \gamma_{\mu} u_{1}\right)}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{2}}-\frac{\left(\bar{u}_{3} \gamma_{\mu} u_{2}\right)\left(\bar{u}_{4} \gamma_{\mu} u_{1}\right)}{\left(\mathbf{p}_{1}-\mathbf{p}_{4}\right)^{2}}\right] \\
& \times\left[\frac{\left(\bar{u}_{4} \gamma_{\nu} u_{2}\right)\left(\bar{u}_{3} \gamma_{\nu} u_{1}\right)}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{2}}-\frac{\left(\bar{u}_{3} \gamma_{\nu} u_{2}\right)\left(\bar{u}_{4} \gamma_{\nu} u_{1}\right)}{\left(\mathbf{p}_{1}-\mathbf{p}_{4}\right)^{2}}\right]^{\dagger} \\
= & 4 \pi^{2} e^{4}\left[\frac{T_{\mathrm{aa}}}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{4}}+\frac{T_{\mathrm{bb}}}{\left(\mathbf{p}_{1}-\mathbf{p}_{4}\right)^{4}}-\frac{T_{\mathrm{ba}}}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{2}\left(\mathbf{p}_{1}-\mathbf{p}_{4}\right)^{2}}\right], \tag{6.82}
\end{align*}
$$

[^45]where
\[

$$
\begin{align*}
T_{\mathrm{aa}} & =\sum_{\text {spins }}\left[\left(\bar{u}_{3} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{3} \gamma_{\nu} u_{1}\right)^{\dagger}\left(\bar{u}_{4} \gamma_{\mu} u_{2}\right)\left(\bar{u}_{4} \gamma_{\nu} u_{2}\right)^{\dagger}\right] \\
& =\operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{1}+m\right) \tilde{\gamma}_{\nu}\left(\not p_{3}+m\right)\right] \operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{2}+m\right) \tilde{\gamma}_{\nu}\left(\not p_{4}+m\right)\right],  \tag{6.83}\\
T_{\mathrm{bb}} & =\sum_{\text {spins }}\left[\left(\bar{u}_{4} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{4} \gamma_{\nu} u_{1}\right)^{\dagger}\left(\bar{u}_{3} \gamma_{\mu} u_{2}\right)\left(\bar{u}_{3} \gamma_{\nu} u_{2}\right)^{\dagger}\right] \\
& =\operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{1}+m\right) \tilde{\gamma}_{\nu}\left(\not p_{4}+m\right)\right] \operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{2}+m\right) \tilde{\gamma}_{\nu}\left(\not p_{3}+m\right)\right],  \tag{6.84}\\
T_{\mathrm{ab}} & =\sum_{\text {spins }}\left[\left(\bar{u}_{3} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{4} \gamma_{\nu} u_{1}\right)^{\dagger}\left(\bar{u}_{4} \gamma_{\mu} u_{2}\right)\left(\bar{u}_{3} \gamma_{\nu} u_{2}\right)^{\dagger}\right] \\
& =\operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{1}+m\right) \tilde{\gamma}_{\nu}\left(\not p_{4}+m\right) \gamma_{\mu}\left(\not p_{2}+m\right) \tilde{\gamma}_{\nu}\left(\not p_{3}+m\right)\right],  \tag{6.85}\\
T_{\mathrm{ba}} & =\sum_{\text {spins }}\left[\left(\bar{u}_{4} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{3} \gamma_{\nu} u_{1}\right)^{\dagger}\left(\bar{u}_{3} \gamma_{\mu} u_{2}\right)\left(\bar{u}_{4} \gamma_{\nu} u_{2}\right)^{\dagger}\right] \\
& =\operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{1}+m\right) \tilde{\gamma}_{\nu}\left(\not p_{3}+m\right) \gamma_{\mu}\left(\not p_{2}+m\right) \tilde{\gamma}_{\nu}\left(\not p_{4}+m\right)\right] . \tag{6.86}
\end{align*}
$$
\]

The symbol $\tilde{\gamma}_{\nu}$ stands for $\tilde{\gamma}_{\nu}=\gamma_{0} \gamma_{\nu}^{\dagger} \gamma_{0}=\gamma_{\nu}$. Expressions (6.83) and (6.84) are obtained by direct application of the formula (2.75) and expressions (6.85) and (6.86) by an extension of the same procedure.

### 6.3.1 Evaluation of traces

To evaluate the traces $T_{\mathrm{aa}}$ and $T_{\mathrm{bb}}$, we use the following trace formula:

$$
\begin{equation*}
\operatorname{Tr}\left[\gamma_{\mu}(\phi+m) \gamma_{\nu}(\not b+m)\right]=4\left[a_{\mu} b_{\nu}+b_{\mu} a_{\nu}-g_{\mu \nu}\left(\mathbf{a} \cdot \mathbf{b}-m^{2}\right)\right] . \tag{6.87}
\end{equation*}
$$

Using Eq. (6.87), we find after simplification

$$
\begin{align*}
T_{\mathrm{aa}}= & 16\left[p_{1 \mu} p_{3 \nu}+p_{3 \mu} p_{1 \nu}-g_{\mu \nu}\left(\mathbf{p}_{1} \cdot \mathbf{p}_{3}-m^{2}\right)\right] \\
& \times\left[p_{2 \mu} p_{4 \nu}+p_{4 \mu} p_{2 \nu}-g_{\mu \nu}\left(\mathbf{p}_{2} \cdot \mathbf{p}_{4}-m^{2}\right)\right] \\
= & 32\left[\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)^{2}+\left(\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right)^{2}+2 m^{2}\left(m^{2}-\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right)\right] \tag{6.88}
\end{align*}
$$

where we have used the relations $g_{\mu \nu} g_{\mu \nu}=4$ and

$$
\mathbf{p}_{1} \cdot \mathbf{p}_{2}=\mathbf{p}_{3} \cdot \mathbf{p}_{4} ; \quad \mathbf{p}_{1} \cdot \mathbf{p}_{3}=\mathbf{p}_{2} \cdot \mathbf{p}_{4} ; \quad \mathbf{p}_{1} \cdot \mathbf{p}_{4}=\mathbf{p}_{2} \cdot \mathbf{p}_{3} ;
$$

which are obtained from the energy-momentum conservation law.

$$
\mathbf{p}_{1}+\mathbf{p}_{2}=\mathbf{p}_{3}+\mathbf{p}_{4} .
$$

Following the same method, $T_{\mathrm{bb}}$ can be evaluated. But, we can write down the result by an inspection of Eq. (6.88) by interchanging $\mathbf{p}_{3} \leftrightarrow \mathbf{p}_{4}$.

$$
\begin{equation*}
T_{\mathrm{bb}}=32\left[\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)^{2}+\left(\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right)^{2}+2 m^{2}\left(m^{2}-\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right)\right], \tag{6.89}
\end{equation*}
$$

The evaluation of $T_{\mathrm{ab}}$ is a little involved. Since the trace of a product of odd number of $\gamma$ matrices vanishes, we need to take only the terms involving even number of $\gamma$ matrices.

$$
\begin{align*}
T_{\mathrm{ab}}= & \operatorname{Tr}\left[\gamma_{\mu} \not p_{1} \gamma_{\nu} \not p_{4} \gamma_{\mu} \not p_{2} \gamma_{\nu} \not p_{3}\right] \\
& +m^{2} \operatorname{Tr}\left[\gamma_{\mu} \not{ }_{1} \gamma_{\nu} \not p_{4} \gamma_{\mu} \gamma_{\nu}+\gamma_{\mu} \not p_{1} \gamma_{\nu} \gamma_{\mu} \not p_{2} \gamma_{\nu}+\gamma_{\mu} \not p_{1} \gamma_{\nu} \gamma_{\mu} \gamma_{\nu} \not p_{3}\right. \\
& \left.+\gamma_{\mu} \gamma_{\nu} \not p_{4} \gamma_{\mu} \not p_{2} \gamma_{\nu}+\gamma_{\mu} \gamma_{\nu} \not p_{4} \gamma_{\mu} \gamma_{\nu} \not p_{3}+\gamma_{\mu} \gamma_{\nu} \gamma_{\mu} \not p_{2} \gamma_{\nu} \not p_{3}\right] \\
& +m^{4} \operatorname{Tr}\left[\gamma_{\mu} \gamma_{\nu} \gamma_{\mu} \gamma_{\nu}\right] . \tag{6.90}
\end{align*}
$$

Using the anticommutation relations of $\gamma$ matrices, we can prove the following formulae involving a product of $\gamma$ matrices of repeated index ${ }^{11}$ (vide Problem 4.1):

$$
\begin{array}{ll}
\gamma_{\mu} \gamma_{\mu}=4 ; & \gamma_{\mu} \gamma_{\nu} \gamma_{\mu}=-2 \gamma_{\nu} ; \\
\gamma_{\mu} \gamma_{\nu} \gamma_{\rho} \gamma_{\mu}=4 g_{\nu \rho} ; & \gamma_{\mu} \gamma_{\nu} \gamma_{\rho} \gamma_{\lambda} \gamma_{\mu}=-2 \gamma_{\lambda} \gamma_{\rho} \gamma_{\nu} . \tag{6.92}
\end{array}
$$

Using the above formulae, we get ${ }^{12}$

$$
\begin{align*}
\underbrace{\gamma_{\mu} \not p_{1} \gamma_{\nu} \not p_{4} \gamma_{\mu}} \not p_{2} \gamma_{\nu} \not p_{3} & =-2 \not p_{4} \underbrace{\gamma_{\nu} \not p_{1} \not p_{2} \gamma_{\nu}} \not p_{3} \\
& =-8\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right) \not p_{4} \not p_{3} ;  \tag{6.93}\\
\operatorname{Tr}\left(\gamma_{\mu} \not p_{1} \gamma_{\nu} \not p_{4} \gamma_{\mu} \not p_{2} \gamma_{\nu} \not p_{3}\right) & =-8\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right) \operatorname{Tr}\left(\not p_{4} \not p_{3}\right) \\
& =-32\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)\left(\mathbf{p}_{3} \cdot \mathbf{p}_{4}\right) ;  \tag{6.94}\\
\underbrace{\gamma_{\mu} \not p_{1} \gamma_{\nu} \not p_{4} \gamma_{\mu}} \gamma_{\nu} & =-2 \not p_{4} \underbrace{\gamma_{\nu} \not p_{1} \gamma_{\nu}}=4 \not p_{4} \not p_{1} ;  \tag{6.95}\\
\operatorname{Tr}\left(\gamma_{\mu} \not p_{1} \gamma_{\nu} \not p_{4} \gamma_{\mu} \gamma_{\nu}\right) & =4 \operatorname{Tr}\left(p_{4} \not p_{1}\right)=16\left(\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right) ;  \tag{6.96}\\
\underbrace{\gamma_{\mu} \gamma_{\nu} \gamma_{\mu}} \gamma_{\nu} & =-2 \gamma_{\nu} \gamma_{\nu}, \quad \text { since } \gamma_{\mu} \gamma_{\nu} \gamma_{\mu}=-2 \gamma_{\nu}  \tag{6.97}\\
\operatorname{Tr}\left(\gamma_{\mu} \gamma_{\nu} \gamma_{\mu} \gamma_{\nu}\right) & =-2 \operatorname{Tr}\left(\gamma_{\mu}\right)^{2}=-8 \operatorname{Tr} I=-32 . \tag{6.98}
\end{align*}
$$

Substituting the above results in Eq. (6.90), we get

$$
\begin{align*}
T_{\mathrm{ab}}= & -32\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)\left(\mathbf{p}_{3} \cdot \mathbf{p}_{4}\right)+16 m^{2}\left\{\mathbf{p}_{1} \cdot \mathbf{p}_{4}+\mathbf{p}_{1} \cdot \mathbf{p}_{2}+\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right. \\
& \left.+\mathbf{p}_{2} \cdot \mathbf{p}_{4}+\mathbf{p}_{3} \cdot \mathbf{p}_{4}+\mathbf{p}_{2} \cdot \mathbf{p}_{3}\right\}-32 m^{4} \\
= & -32\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)^{2}+32 m^{2}\left\{\mathbf{p}_{1} \cdot \mathbf{p}_{4}+\mathbf{p}_{1} \cdot \mathbf{p}_{2}+\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right\}-32 m^{4} \tag{6.99}
\end{align*}
$$

[^46]Equation (6.99) remains the same under the exchange $\mathbf{p}_{3} \leftrightarrow \mathbf{p}_{4}$. So, the trace $T_{\mathrm{ba}}=T_{\mathrm{ab}}$. Hence

$$
\begin{align*}
T_{\mathrm{ab}}+T_{\mathrm{ba}} & =64\left[-\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)^{2}+m^{2}\left\{\mathbf{p}_{1} \cdot \mathbf{p}_{4}+\mathbf{p}_{1} \cdot \mathbf{p}_{2}+\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right\}-m^{4}\right] \\
& =64\left[-\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)^{2}+2 m^{2}\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)\right] . \tag{6.100}
\end{align*}
$$

The last step is obtained by invoking the energy-momentum conservation.

$$
\begin{align*}
\mathbf{p}_{1} \cdot \mathbf{p}_{4} & =\mathbf{p}_{1} \cdot\left(\mathbf{p}_{1}+\mathbf{p}_{2}-\mathbf{p}_{3}\right) \\
& =m^{2}+\mathbf{p}_{1} \cdot \mathbf{p}_{2}-\mathbf{p}_{1} \cdot \mathbf{p}_{3} . \tag{6.101}
\end{align*}
$$

The denominators in Eq. (6.82) can be expressed in terms of scalar products of four vectors.

$$
\begin{align*}
&\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{2}=\mathbf{p}_{1}^{2}+\mathbf{p}_{3}^{2}-2 \mathbf{p}_{1} \cdot \mathbf{p}_{3}=2\left(m^{2}-\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right) .  \tag{6.102}\\
&\left(\mathbf{p}_{1}-\mathbf{p}_{4}\right)^{2}=\mathbf{p}_{1}^{2}+\mathbf{p}_{4}^{2}-2 \mathbf{p}_{1} \cdot \mathbf{p}_{4}=2\left(m^{2}-\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right) .  \tag{6.103}\\
&|\overline{\mathcal{M}}|^{2}= 32 \pi^{2} e^{4}\left[\frac{\left\{\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)^{2}+\left(\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right)^{2}+2 m^{2}\left(m^{2}-\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right)\right\}}{\left(m^{2}-\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right)^{2}}\right. \\
&+\frac{\left\{\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)^{2}+\left(\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right)^{2}+2 m^{2}\left(m^{2}-\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right)\right\}}{\left(m^{2}-\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right)^{2}} \\
&\left.-\frac{2\left\{-\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)^{2}+2 m^{2}\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)\right\}}{\left(m^{2}-\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right)\left(m^{2}-\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right)}\right] \tag{6.104}
\end{align*}
$$

Equation (6.104) is valid in all frames of reference, both in laboratory and centre of momentum frames.

### 6.3.2 In the centre of momentum frame

Let us choose the centre of momentum (c.m.) frame for evaluation of the cross section. It describes the elastic scattering of two electrons as shown in Fig. 6.5 with the three-momenta of incident electrons equal and opposite $\left(\boldsymbol{p}_{2}=-\boldsymbol{p}_{1}\right)$. Similarly, the scattered electrons have also their three-momenta equal and opposite $\left(\boldsymbol{p}_{4}=-\boldsymbol{p}_{3}\right)$. Also the energies of the incident and scattered electrons are equal.

$$
\begin{equation*}
\left|\boldsymbol{p}_{1}\right|=\left|\boldsymbol{p}_{2}\right|=\left|\boldsymbol{p}_{3}\right|=\left|\boldsymbol{p}_{4}\right|=p \quad \text { and } \quad E_{1}=E_{2}=E_{3}=E_{4}=E . \tag{6.105}
\end{equation*}
$$



Figure 6.5: Electron-electron scattering in c.m. frame. $\boldsymbol{p}_{1}$ and $\boldsymbol{p}_{2}$ denote the three-momenta before collision and $\boldsymbol{p}_{3}$ and $\boldsymbol{p}_{4}$ denote the three-momenta after collision, $\theta$ being the angle of scattering.

Now we can write the scalar product of the four-vectors and the other quantities in terms of $E, p$ and the angle of scattering $\theta$.

$$
\begin{align*}
\mathbf{p}_{1} \cdot \mathbf{p}_{2} & =E_{1} E_{2}-\boldsymbol{p}_{1} \cdot \boldsymbol{p}_{2}=E^{2}+p^{2}  \tag{6.106}\\
\mathbf{p}_{1} \cdot \mathbf{p}_{3} & =E_{1} E_{3}-\boldsymbol{p}_{1} \cdot \boldsymbol{p}_{3}=E^{2}-p^{2} \cos \theta  \tag{6.107}\\
\mathbf{p}_{1} \cdot \mathbf{p}_{4} & =E_{1} E_{4}-\boldsymbol{p}_{1} \cdot \boldsymbol{p}_{4}=E^{2}+p^{2} \cos \theta  \tag{6.108}\\
\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{2} & =2 m^{2}\left(1-\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right)=-2 p^{2}(1-\cos \theta)  \tag{6.109}\\
\left(\mathbf{p}_{1}-\mathbf{p}_{4}\right)^{2} & =2 m^{2}\left(1-\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right)=-2 p^{2}(1+\cos \theta) \tag{6.110}
\end{align*}
$$

Substituting the results (6.106) - (6.110) into Eqs. (6.104), we obtain an expression for $|\overline{\mathcal{M}}|^{2}$ defined in Eq. (6.82).

$$
\begin{align*}
|\overline{\mathcal{M}}|^{2}= & \frac{32 \pi^{2} e^{4}}{p^{4}}\left[\frac{\left\{\left(E^{2}+p^{2}\right)^{2}+\left(E^{2}+p^{2} \cos \theta\right)^{2}-2 m^{2} p^{2}(1-\cos \theta)\right\}}{(1-\cos \theta)^{2}}\right. \\
& +\frac{\left\{\left(E^{2}+p^{2}\right)^{2}+\left(E^{2}-p^{2} \cos \theta\right)^{2}-2 m^{2} p^{2}(1+\cos \theta)\right\}}{(1+\cos \theta)^{2}} \\
& \left.+\frac{2\left\{\left(E^{2}+p^{2}\right)^{2}-2 m^{2}\left(E^{2}+p^{2}\right)\right\}}{\sin ^{2} \theta}\right] \tag{6.111}
\end{align*}
$$

## Differential cross section in c.m. frame

The cross section is given by

$$
\begin{equation*}
d \sigma=\frac{2 \pi}{v}(\Pi N)^{-1}|\overline{\mathcal{M}}|^{2} \rho_{f}, \tag{6.112}
\end{equation*}
$$

where $(\Pi N)^{-1}$ denotes the normalization factor

$$
(\Pi N)^{-1}=\frac{1}{2 E_{1} 2 E_{2} 2 E_{3} 2 E_{4}}=\left(\frac{1}{2 E}\right)^{4}
$$

and $\rho_{f}$ is the density of final states for the two-particle system (vide Appendix).

$$
\begin{align*}
\rho_{f} & =\frac{1}{(2 \pi)^{3}} E_{3} E_{4} \frac{p_{3}^{3} d \Omega}{\left(E_{3}+E_{4}\right) p_{3}^{2}-E_{3}\left(\boldsymbol{P} \cdot \boldsymbol{p}_{3}\right)} \\
& =\frac{1}{(2 \pi)^{3}} E^{2} \frac{p^{3} d \Omega}{2 E p^{2}}=\frac{E p}{2(2 \pi)^{3}} d \Omega, \tag{6.113}
\end{align*}
$$

where $\boldsymbol{P}=\boldsymbol{p}_{3}+\boldsymbol{p}_{4}$ denotes the total momentum of the two-particle system which is zero in the c.m. frame. The incident flux depends upon the relative velocity $v$ of the two colliding particles and in the c.m. system, it is equal to twice the velocity of either of the colliding partners which are of equal mass.

$$
v=2 \frac{p}{E}
$$

Collecting all these results, we arrive at an expression for the differential cross section.

$$
\begin{align*}
\frac{d \sigma}{d \Omega} & =\frac{2 \pi}{2(p / E)}\left(\frac{1}{2 E}\right)^{4}|\overline{\mathcal{M}}|^{2} \frac{E p}{2(2 \pi)^{3}} \\
& =\frac{1}{64 E^{2}} \frac{1}{(2 \pi)^{2}}|\overline{\mathcal{M}}|^{2} . \tag{6.114}
\end{align*}
$$

Equation (6.114) together with Eq. (6.111) gives the electron-electron scattering cross section in c.m. frame. The electron-electron scattering is sometimes known as $\mathrm{M} \phi 1 \mathrm{ller}$ scattering since it was $\mathrm{M} \phi 1 \mathrm{ler}{ }^{13}$ who first deduced the electron-electron scattering cross section in the laboratory frame using relativistic quantum mechanics.

### 6.4 Electron-positron scattering

Electron-positron scattering is often referred to as Bhabha scattering since it was he who first deduced the relativistic cross section for this process ${ }^{14}$. Electron-positron scattering is similar to electron-electron scattering but the arrows representing the positron are reversed, indicating that the positron travels backward in time according to Feynman's positron theory. In this case also, there are two Feynman diagrams, one indicating an exchange of a photon between electron and positron and the other

[^47]

Figure 6.6: Electron-positron scattering. Diagram (a) represents the exchange of one photon between electron and positron and Diagram (b) denotes the $e^{+} e^{-}$ pair annihilation and subsequent creation. Both the diagrams contribute to the process.
indicating electron-positron pair annihilation into a photon which subsequently creates an electron-positron pair. Both the diagrams contribute to the Feynman amplitude for the process. Using the short-hand notation $u_{1}, \bar{v}_{2}, \bar{u}_{3}, v_{4}$ for the spinors $u\left(\mathbf{p}_{1}\right), \bar{v}\left(\mathbf{p}_{2}\right), \bar{u}\left(\mathbf{p}_{3}\right), v\left(\mathbf{p}_{4}\right)$, we write down the matrix element for the electron-positron scattering.

$$
\begin{equation*}
\mathcal{M}=i 4 \pi e^{2}\left[\frac{\left(\bar{u}_{3} \gamma_{\mu} u_{1}\right)\left(\bar{v}_{2} \gamma_{\mu} v_{4}\right)}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{2}}-\frac{\left(\bar{v}_{2} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{3} \gamma_{\mu} v_{4}\right)}{\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{2}}\right], \tag{6.115}
\end{equation*}
$$

The relative negative sign between the two amplitudes arises from fermion exchange - the initial $e^{+}$(negative energy electron) is exchanged with the final $e^{-}$. To calculate the cross section for unpolarized particles, we need to find the absolute square of the above amplitude (6.115) by summing over the final spin states and averaging over the initial spin states.

$$
\begin{align*}
|\overline{\mathcal{M}}|^{2}= & \frac{1}{4} \mathcal{M} \mathcal{M}^{\dagger} \\
= & 4 \pi^{2} e^{4} \sum_{\text {spins }}
\end{aligned} \begin{aligned}
& {\left[\frac{\left\{\left(\bar{u}_{3} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{1} \gamma_{\nu} u_{3}\right)\right\}\left\{\left(\bar{v}_{2} \gamma_{\mu} v_{4}\right)\left(\bar{v}_{4} \gamma_{\nu} v_{2}\right)\right\}}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{4}}\right.} \\
& +\frac{\left\{\left(\bar{v}_{2} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{1} \gamma_{\nu} v_{2}\right)\right\}\left\{\left(\bar{u}_{3} \gamma_{\mu} v_{4}\right)\left(\bar{v}_{4} \gamma_{\nu} u_{3}\right)\right\}}{\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{4}} \\
& -\frac{\left\{\left(\bar{u}_{3} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{1} \gamma_{\nu} v_{2}\right)\left(\bar{v}_{2} \gamma_{\mu} v_{4}\right)\left(\bar{v}_{4} \gamma_{\nu} u_{3}\right)\right\}}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{2}\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{2}} \\
& \left.\quad-\frac{\left\{\left(\bar{v}_{2} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{1} \gamma_{\nu} u_{3}\right)\left(\bar{u}_{3} \gamma_{\mu} v_{4}\right)\left(\bar{v}_{4} \gamma_{\nu} v_{2}\right)\right\}}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{2}\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{2}}\right] \\
= & 4 \pi^{2} e^{4}\left[\frac{T_{\mathrm{aa}}}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{4}}+\frac{T_{\mathrm{bb}}}{\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{4}}-\frac{T_{\mathrm{ab}}+T_{\mathrm{ba}}}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{2}\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{2}}\right], \tag{6.116}
\end{align*}
$$

where

$$
\begin{align*}
T_{\mathrm{aa}} & =\sum_{\text {spins }}\left[\left\{\left(\bar{u}_{3} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{1} \gamma_{\nu} u_{3}\right)\right\}\left\{\left(\bar{v}_{2} \gamma_{\mu} v_{4}\right)\left(\bar{v}_{4} \gamma_{\nu} v_{2}\right)\right\}\right] \\
& =\operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{1}+m\right) \gamma_{\nu}\left(\not p_{3}+m\right)\right] \operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{4}-m\right) \gamma_{\nu}\left(\not p_{2}-m\right)\right] ;  \tag{6.117}\\
T_{\mathrm{bb}} & =\sum_{\text {spins }}\left[\left\{\left(\bar{v}_{2} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{1} \gamma_{\nu} v_{2}\right)\right\}\left\{\left(\bar{u}_{3} \gamma_{\mu} v_{4}\right)\left(\bar{v}_{4} \gamma_{\nu} u_{3}\right)\right\}\right] \\
& =\operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{1}+m\right) \gamma_{\nu}\left(\not p_{2}-m\right)\right] \operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{4}-m\right) \gamma_{\nu}\left(\not \not{ }_{3}+m\right)\right] ;  \tag{6.118}\\
T_{\mathrm{ab}} & =\sum_{\text {spins }}\left[\left(\bar{u}_{3} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{1} \gamma_{\nu} v_{2}\right)\left(\bar{v}_{2} \gamma_{\mu} v_{4}\right)\left(\bar{v}_{4} \gamma_{\nu} u_{3}\right)\right] \\
& =\operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{1}+m\right) \gamma_{\nu}\left(\not p_{2}-m\right) \gamma_{\mu}\left(\not p_{4}-m\right) \gamma_{\nu}\left(\not p_{3}+m\right)\right] ;  \tag{6.119}\\
T_{\mathrm{ba}} & =\sum_{\text {spins }}\left[\left(\bar{v}_{2} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{1} \gamma_{\nu} u_{3}\right)\left(\bar{u}_{3} \gamma_{\mu} v_{4}\right)\left(\bar{v}_{4} \gamma_{\nu} v_{2}\right)\right] \\
& =\operatorname{Tr}\left[\gamma_{\mu}\left(\not \not{ }_{1}+m\right) \gamma_{\nu}\left(\not p_{3}+m\right) \gamma_{\mu}\left(\not p_{4}-m\right) \gamma_{\nu}\left(\not p_{2}-m\right)\right] . \tag{6.120}
\end{align*}
$$

These traces can be calculated by following the same procedure used in the previous section on electron-electron scattering. By an inspection, it is found that the traces can be written down directly from Eqs. (6.88), (6.89), (6.99) and (6.100) by interchanging $\mathbf{p}_{2} \leftrightarrow-\mathbf{p}_{4}$.

$$
\begin{align*}
& T_{\mathrm{aa}}=32\left[\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)^{2}+\left(\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right)^{2}+2 m^{2}\left(m^{2}-\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right)\right],  \tag{6.121}\\
& T_{\mathrm{bb}}=32\left[\left(\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right)^{2}+\left(\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right)^{2}+2 m^{2}\left(m^{2}+\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)\right],  \tag{6.122}\\
& T_{\mathrm{ab}}=-32\left(\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right)^{2}-32 m^{2}\left\{\mathbf{p}_{1} \cdot \mathbf{p}_{2}+\mathbf{p}_{1} \cdot \mathbf{p}_{4}-\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right\}-32 m^{4} .(6.122)  \tag{6.123}\\
& T_{\mathrm{ba}}=-32\left(\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right)^{2}-32 m^{2}\left\{\mathbf{p}_{1} \cdot \mathbf{p}_{2}+\mathbf{p}_{1} \cdot \mathbf{p}_{4}-\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right\}-32 m^{4} .(6.124) \tag{6.124}
\end{align*}
$$

## Differential cross section in c.m. frame

Till now, we have not specified the frames of reference and so the expressions (6.121) - (6.124) are valid in all frames of reference. Now, we shall express the scalar products of four-vectors in c.m. frame as given by Eqs. (6.106) - (6.110) and substitute them in Eq. (6.116) to get $|\overline{\mathcal{M}}|^{2}$.

$$
\begin{align*}
|\overline{\mathcal{M}}|^{2}=32 \pi^{2} e^{4} & {\left[\frac{\left(E^{2}+p^{2}\right)^{2}+\left(E^{2}+p^{2} \cos \theta\right)^{2}-2 m^{2} p^{2}(1-\cos \theta)}{p^{4}(1-\cos \theta)^{2}}\right.} \\
& +\frac{E^{4}+p^{4} \cos ^{2} \theta+2 m^{2} E^{2}}{2 E^{4}} \\
& \left.-\frac{\left(E^{2}+p^{2} \cos \theta\right)^{2}+2 m^{2}\left(E^{2}+p^{2} \cos \theta\right)}{E^{2} p^{2}(1-\cos \theta)}\right] . \tag{6.125}
\end{align*}
$$

By substituting expression (6.125) in Eq. (6.114), we obtain the differential cross section for electron - positron scattering.

$$
\begin{align*}
\frac{d \sigma}{d \Omega}=\frac{e^{4}}{8 E^{2}}[ & \frac{\left(E^{2}+p^{2}\right)^{2}+\left(E^{2}+p^{2} \cos \theta\right)^{2}-2 m^{2} p^{2}(1-\cos \theta)}{p^{4}(1-\cos \theta)^{2}} \\
& +\frac{E^{4}+p^{4} \cos ^{2} \theta+2 m^{2} E^{2}}{2 E^{4}} \\
& \left.-\frac{\left(E^{2}+p^{2} \cos \theta\right)^{2}+2 m^{2}\left(E^{2}+p^{2} \cos \theta\right)}{E^{2} p^{2}(1-\cos \theta)}\right] \tag{6.126}
\end{align*}
$$

where $e^{2}$ is the fine structure constant $\alpha$. In the high energy limit ( $E \gg$ $m$ ), Eq. (6.126) reduces to

$$
\begin{equation*}
\frac{d \sigma}{d \Omega}=\frac{e^{4}}{8 E^{2}}\left[\frac{1+\cos ^{4}(\theta / 2)}{\sin ^{4}(\theta / 2)}+\frac{1+\cos ^{2} \theta}{2}-\frac{2 \cos ^{4}(\theta / 2)}{\sin ^{2}(\theta / 2)}\right] . \tag{6.127}
\end{equation*}
$$

### 6.5 Electron-positron pair annihilation into two photons

We have observed earlier that the first order process of electron-positron pair annihilation into a single photon is unphysical and so we need to consider necessarily a second order process of pair annihilation into two photons. This is depicted by Feynman diagrams in Fig. 6.7. The diagram (b) represents the exchange diagram. The Feynman amplitudes for


Figure 6.7: Electron-positron pair annihilation into two photons. The internal line connecting the two vertices denotes a virtual particle.
diagrams (a) and (b) can be written as

$$
\begin{gather*}
\mathcal{M}_{\mathrm{a}}=-i 4 \pi e^{2}\left(\bar{v}_{2} \not \xi_{2} \frac{1}{\not p_{1}-\not \phi_{1}-m} \xi_{1} u_{1}\right),  \tag{6.128}\\
\mathcal{M}_{\mathrm{b}}=-i 4 \pi e^{2}\left(\bar{v}_{2} \not_{1} \frac{1}{\not p_{1}-\not \phi_{2}-m} \phi_{2} u_{1}\right),  \tag{6.129}\\
\mathcal{M}=\mathcal{M}_{\mathrm{a}}+\mathcal{M}_{\mathrm{b}}=-i 4 \pi e^{2}\left\{\bar{v}_{2}\left(\mathcal{O}_{\mathrm{a}}+\mathcal{O}_{\mathrm{b}}\right) u_{1}\right\}, \tag{6.130}
\end{gather*}
$$

where

$$
\begin{align*}
\mathcal{O}_{\mathrm{a}} & =\not \xi_{2} \frac{1}{\not p_{1}-\not k_{1}-m} \xi_{1}=\not \phi_{2} \frac{\not p_{1}-\not \phi_{1}+m}{\left(\not p_{1}-\not 1_{1}\right)^{2}-m^{2}} \xi_{1}  \tag{6.131}\\
\mathcal{O}_{\mathrm{b}} & =\not \phi_{1} \frac{1}{\not p_{1}-\not \phi_{2}-m} \xi_{2}=\not \phi_{1} \frac{\not p_{1}-\not \phi_{2}+m}{\left(\not p_{1}-\not k_{2}\right)^{2}-m^{2}} \phi_{2} . \tag{6.132}
\end{align*}
$$

Let us now choose a coordinate system in which the electron is at rest. By this choice, the matrix element calculation is greatly simplified. The denominators in Eqs. (6.131) and (6.132) simplify to

$$
\begin{align*}
\left(\not p_{1}-\not \not p_{1}\right)^{2}-m^{2} & =\not p_{1}^{2}+\not \psi_{1}^{2}-2 \not p_{1} \not k_{1}-m^{2} \\
& =\mathbf{p}_{1}^{2}+\mathbf{k}_{1}^{2}-2 \mathbf{p}_{1} \cdot \mathbf{k}_{1}-m^{2}=-2 m \omega_{1}  \tag{6.133}\\
\left(\not p_{1}-\not \not p_{2}\right)^{2}-m^{2} & =\not p_{1}^{2}+\not 2_{2}^{2}-2 \not p_{1} \not \phi_{2}-m^{2} \\
& =\mathbf{p}_{1}^{2}+\mathbf{k}_{2}^{2}-2 \mathbf{p}_{1} \cdot \mathbf{k}_{2}-m^{2}=-2 m \omega_{2} \tag{6.134}
\end{align*}
$$

The numerators can also be similarly simplified. The first term in the numerators of Eqs. (6.131) and (6.132) taken between the electron and positron spinors

$$
\begin{align*}
& \bar{v}_{2} \xi_{2} \underbrace{\not p_{1} \xi_{1}} u_{1}=\bar{v}_{2} \xi_{2}\left(-\phi_{1} \not b_{1}+2 \mathbf{p}_{1} \cdot \varepsilon_{1}\right) u_{1} \\
& =-\bar{v}_{2} \xi_{2} \xi_{1} \not 1_{1} u_{1}=-m \bar{v}_{2} \xi_{2} \xi_{1} u_{1},  \tag{6.135}\\
& \bar{v}_{2} \not_{1} p_{1} \not_{2} u_{1}=\bar{v}_{2} \phi_{1}\left(-\not \phi_{2} \not p_{1}+2 \mathbf{p}_{1} \cdot \varepsilon_{2}\right) u_{1} \\
& =-\bar{v}_{2} \xi_{1} \xi_{2} \not p_{1} u_{1}=-m \bar{v}_{2} \xi_{1} \xi_{2} u_{1} \text {, } \tag{6.136}
\end{align*}
$$

cancels with the third term. Thus we obtain

$$
\begin{align*}
& \bar{v}_{2} \mathcal{O}_{\mathrm{a}} u_{1}=-\frac{\bar{v}_{2} \xi_{2} \not k_{1} \xi_{1} u_{1}}{2 m \omega_{1}}=\frac{\bar{v}_{2} \xi_{2} \xi_{1} \not k_{1} u_{1}}{2 m \omega_{1}} ;  \tag{6.137}\\
& \bar{v}_{2} \mathcal{O}_{\mathrm{b}} u_{1}=-\frac{\bar{v}_{2} \xi_{1} \not k_{2} \xi_{2} u_{1}}{2 m \omega_{2}}=\frac{\bar{v}_{2} \phi_{1} \xi_{2} \not k_{2} u_{1}}{2 m \omega_{2}} . \tag{6.138}
\end{align*}
$$

## Transition probability

If we are not interested in the spin states of the electron and positron, then an average has to be taken over the spin states in calculating the transition probability.

$$
\begin{equation*}
|\overline{\mathcal{M}}|^{2}=16 \pi^{2} e^{4} \sum^{-}\left|\bar{v}_{2} \mathcal{O} u_{1}\right|^{2} \tag{6.139}
\end{equation*}
$$

where

$$
\begin{align*}
\mathcal{O} & =\mathcal{O}_{\mathrm{a}}+\mathcal{O}_{\mathrm{b}} \\
& =\frac{\xi_{2} \xi_{1} k_{1}}{2 m \omega_{1}}+\frac{\xi_{1} \xi_{2} \not \phi_{2}}{2 m \omega_{1}} . \tag{6.140}
\end{align*}
$$

Using the trace techniques, we get

$$
\begin{align*}
\bar{\sum}\left|\bar{v}_{2} \mathcal{O} u_{1}\right|^{2} & =\frac{1}{4} \operatorname{Tr}\left[\mathcal{O}\left(\not p_{1}+m\right) \tilde{\mathcal{O}}\left(\not p_{2}-m\right)\right] \\
& =\frac{1}{4} \operatorname{Tr}\left[\left(\mathcal{O}_{\mathrm{a}}+\mathcal{O}_{\mathrm{b}}\right)\left(\not p_{1}+m\right)\left(\tilde{\mathcal{O}}_{\mathrm{a}}+\tilde{\mathcal{O}}_{\mathrm{b}}\right)\left(\not p_{2}-m\right)\right] \\
& =\frac{1}{4}\left[T_{\mathrm{aa}}+T_{\mathrm{bb}}+T_{\mathrm{ab}}+T_{\mathrm{ba}}\right], \tag{6.141}
\end{align*}
$$

where

$$
\begin{align*}
& T_{\mathrm{aa}}=\operatorname{Tr}\left[\mathcal{O}_{\mathrm{a}}\left(\not p_{1}+m\right) \tilde{\mathcal{O}}_{\mathrm{a}}\left(\not p_{2}-m\right)\right],  \tag{6.142}\\
& T_{\mathrm{bb}}=\operatorname{Tr}\left[\mathcal{O}_{\mathrm{b}}\left(\not p_{1}+m\right) \tilde{\mathcal{O}}_{\mathrm{b}}\left(\not p_{2}-m\right)\right],  \tag{6.143}\\
& T_{\mathrm{ab}}=\operatorname{Tr}\left[\mathcal{O}_{\mathrm{a}}\left(\not p_{1}+m\right) \tilde{\mathcal{O}}_{\mathrm{b}}\left(\not p_{2}-m\right)\right],  \tag{6.144}\\
& T_{\mathrm{ba}}=\operatorname{Tr}\left[\mathcal{O}_{\mathrm{b}}\left(\not p_{1}+m\right) \tilde{\mathcal{O}}_{\mathrm{a}}\left(\not p_{2}-m\right)\right], \tag{6.145}
\end{align*}
$$

with

$$
\begin{align*}
& \mathcal{O}_{\mathrm{a}}=\frac{\xi_{2} \xi_{1} \not \phi_{1}}{2 m \omega_{1}}, \quad \mathcal{O}_{\mathrm{b}}=\frac{\xi_{1} \xi_{2} \not \phi_{2}}{2 m \omega_{1}},  \tag{6.146}\\
& \tilde{\mathcal{O}}_{\mathrm{a}}=\frac{\not{ }_{1} \phi_{1} \phi_{1} \xi_{2}}{2 m \omega_{1}}, \quad \tilde{\mathcal{O}}_{\mathrm{b}}=\frac{\not \not \xi_{2} \phi_{2} \phi_{1}}{2 m \omega_{2}} . \tag{6.147}
\end{align*}
$$

The traces are similar to those encountered in the study of Compton scattering and they can be evaluated following the same procedure. From the results obtained in Compton scattering, one can write down the traces
for pair annihilation into two photons by substituting $\mathbf{k}_{1} \rightarrow-\mathbf{k}_{1}$ and $\omega_{1} \rightarrow-\omega_{1}$.

$$
\begin{align*}
& T_{\mathrm{aa}}=-\frac{2}{m \omega_{1}}\left\{m \omega_{2}-2\left(-\mathbf{k}_{1} \cdot \varepsilon_{2}\right)^{2}\right\} .  \tag{6.148}\\
& T_{\mathrm{bb}}=\frac{2}{m \omega_{2}}\left\{-m \omega_{1}-2\left(\mathbf{k}_{2} \cdot \varepsilon_{1}\right)^{2}\right\} .  \tag{6.149}\\
& T_{\mathrm{ab}}=2\left\{2\left(\varepsilon_{1} \cdot \varepsilon_{2}\right)^{2}-1\right\}-\frac{2}{m \omega_{1}}\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)^{2}+\frac{2}{m \omega_{2}}\left(\mathbf{k}_{2} \cdot \varepsilon_{1}\right)^{2} .  \tag{6.150}\\
& T_{\mathrm{ba}}=T_{\mathrm{ab}} . \tag{6.151}
\end{align*}
$$

## Density of final states

For the final states of two particles $a$ and $b$, the density of states $\rho_{f}$ is given by (vide Appendix B)

$$
\rho_{f}=(2 \pi)^{-3} E_{a} E_{b} \frac{p_{a}^{3} d \Omega_{a}}{E p_{a}^{2}-E_{a}\left(\boldsymbol{p} \cdot \boldsymbol{p}_{a}\right)},
$$

where $E$ denotes the total energy and $\boldsymbol{p}$ denotes the total momentum. In the present case for diagram (a),

$$
E_{a}=\omega_{1} ; \quad E_{b}=\omega_{2} ; \quad E=\omega_{1}+\omega_{2} ; \quad \boldsymbol{p}=\boldsymbol{k}_{1}+\boldsymbol{k}_{2} ; \quad \boldsymbol{p}_{a}=\boldsymbol{k}_{1} ; \quad\left|\boldsymbol{p}_{a}\right|=\omega_{1} .
$$

Invoking energy-momentum conservation, we have

$$
\begin{align*}
\mathbf{p}_{1}+\mathbf{p}_{2} & =\mathbf{k}_{1}+\mathbf{k}_{2} ; \\
\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{2} & =\left(\mathbf{k}_{1}+\mathbf{k}_{2}\right)^{2} ; \\
\mathbf{p}_{1}^{2}+\mathbf{p}_{2}^{2}+2 \mathbf{p}_{1} \cdot \mathbf{p}_{2} & =\mathbf{k}^{2}+\mathbf{k}_{2}^{2}+\mathbf{k}_{1} \cdot \mathbf{k}_{2} ; \\
2 m^{2}+2 m E_{2} & =2 \omega_{1} \omega_{2}(1-\cos \theta) . \tag{6.152}
\end{align*}
$$

Using the above results, we get

$$
\begin{align*}
\rho_{f} & =\frac{\omega_{1} \omega_{2}}{(2 \pi)^{3}} \frac{\omega_{1}^{3} d \Omega_{1}}{\left(\omega_{1}+\omega_{2}\right) \omega_{1}^{2}-\omega_{1}\left(\omega_{1}^{2}+\boldsymbol{k}_{1} \cdot \boldsymbol{k}_{2}\right)} \\
& =\frac{\omega_{1}^{3} \omega_{2} d \Omega_{1}}{(2 \pi)^{3}\left(\omega_{1} \omega_{2}-\boldsymbol{k}_{1} \cdot \boldsymbol{k}_{2}\right)} . \tag{6.153}
\end{align*}
$$

Using the relation $\boldsymbol{k}_{1} \cdot \boldsymbol{k}_{2}=\omega_{1} \omega_{2} \cos \theta$ and Eq. (6.152), we get

$$
\begin{equation*}
\rho_{f}=\frac{\omega_{1}^{3} \omega_{2} d \Omega_{1}}{(2 \pi)^{3} m\left(m+E_{2}\right)} \tag{6.154}
\end{equation*}
$$

If $v$ is the velocity of the incident positron, then the differential cross section is given by

$$
\begin{align*}
d \sigma & =\frac{2 \pi}{v}\left(\Pi_{N}\right)|\overline{\mathcal{M}}|^{2} \rho_{f} \\
& =\frac{2 \pi}{\left(p_{2} / E_{2}\right)} \frac{1}{2 E_{1} 2 E_{2} 2 \omega_{1} 2 \omega_{2}}|\overline{\mathcal{M}}|^{2} \frac{\omega_{1}^{3} \omega_{2} d \Omega_{1}}{(2 \pi)^{3} m\left(m+E_{2}\right)} \\
& =\frac{\omega_{1}^{2}}{64 \pi^{2} m^{2} p_{2}\left(m+E_{2}\right)}|\overline{\mathcal{M}}|^{2} d \Omega_{1} . \tag{6.155}
\end{align*}
$$

### 6.6 Bremsstrahlung

Let us consider a moving electron deflected by the Coulomb potential of a nucleus. This will cause an emission of radiation according to the classical theory. This process is known as Bremsstrahlung. We shall calculate the transition probability for the electron to go from one state to another with the emission of a photon in the field of a nucleus. Interaction with an external field is essential for the simultaneous conservation of energy and momentum for the process. In other words, an electron cannot make a transition from one state to another emitting a photon while travelling in a vacuum.

Assuming the Born approximation in which the Coulomb potential acts only once on the electron, there are two modes by which the electron can emit a photon. One mode is that the electron is first scattered by the Coulomb potential and then emits a photon to reach the final state. The other mode is that the electron emits a photon first and then interacts with the Coulomb field. Both these modes are indistinguishable and hence contribute to the total amplitude of transition for the electron from the initial to the final state. The Feynman diagrams for both the modes are depicted in Fig. 6.8.

By an inspection of the two Feynman diagrams (a) and (b), one can write down the total matrix element for the process.

$$
\left.\begin{array}{rl}
\mathcal{M}=-i\left(4 \pi e^{2}\right)^{1 / 2} & {\left[\bar{u}\left(\mathbf{p}_{2}\right) \notin(\mathbf{k}) \frac{1}{\not p_{1}+\not q-m} \phi(\mathbf{q}) u\left(\mathbf{p}_{1}\right)\right.} \\
& +\bar{u}\left(\mathbf{p}_{2}\right) \phi(\mathbf{q}) \frac{1}{\not p_{2}-\not q-m} \tag{6.156}
\end{array}(\mathbf{k}) u\left(\mathbf{p}_{1}\right)\right], ~ \$
$$

with

$$
\phi(\mathbf{q})=\frac{4 \pi z e^{2}}{q^{2}} \delta\left(\mathbf{q}_{0}\right) \gamma_{0} \quad \text { and } \quad \notin=\sum_{\mu} \epsilon_{\mu} \gamma_{\mu} .
$$



Figure 6.8: Feynman diagrams representing Bremsstrahlung process.

The electron propagators can be rationalized to yield

$$
\begin{align*}
& \frac{1}{\not p_{1}+\not q-m}=\frac{\not p_{1}+\not q+m}{\left(\not p_{1}+\not q\right)^{2}-m^{2}}=\frac{\not p_{1}+\not q+m}{-2 \boldsymbol{p}_{1} \cdot \boldsymbol{q}-q^{2}},  \tag{6.157}\\
& \frac{1}{\not p_{2}-\not q-m}=\frac{\not p_{2}-\not q+m}{\left(\not p_{2}-\not q\right)^{2}-m^{2}}=\frac{\not{ }_{2}-\not q+m}{2 \boldsymbol{p}_{2} \cdot \boldsymbol{q}-q^{2}}, \tag{6.158}
\end{align*}
$$

where we have used the condition that in Coulomb scattering, the energy transfer to the electron $q_{0}=0$ and the four-momentum square for the electron $\mathbf{p}_{1}^{2}=m^{2}$ and $\mathbf{p}_{2}^{2}=m^{2}$. Using the schematic diagram for the Bremsstrahlung process, the denominators in Eqs. (6.157) and (6.158) can be simplified to yield

$$
\begin{align*}
& D_{1}=-2 \boldsymbol{p}_{1} \cdot \boldsymbol{q}-q^{2}=2 k\left(E_{2}-p_{2} \cos \theta_{2}\right)  \tag{6.159}\\
& D_{2}=2 \boldsymbol{p}_{2} \cdot \boldsymbol{q}-q^{2}=-2 k\left(E_{1}-p_{1} \cos \theta_{1}\right) \tag{6.160}
\end{align*}
$$

where $E_{1}, E_{2}$ denote the energies of the incident and outgoing electron and $\theta_{1}, \theta_{2}$ denote the angles of the emitted photon and outgoing electron as shown in the schematic diagram 6.9.

Summing over the two spin states of the outgoing electron and the two polarization states of the emitted photon and averaging over the initial spin states of the electron, we can evaluate the differential cross section for the process using the Fermi's golden rule.

$$
\begin{equation*}
d \sigma=\frac{2 \pi}{v_{1}} \frac{1}{2 E_{1} 2 E_{2} 2 \omega}|\overline{\mathcal{M}}|^{2} \rho_{f}, \tag{6.161}
\end{equation*}
$$



Figure 6.9: Schematic diagram representing the Bremsstrahlung process. The emitted photon makes an angle $\theta_{1}$ with the incident electron and $\theta_{2}$ denotes the direction of the outgoing electron. The two planes, one containing incoming electron and outgoing photon and the other containing incoming electron and outgoing electron need not be the same and the angle between them is taken as $\phi$. The two planes are not shown explicitly in figure.
where $v_{1}$ denotes the energy of the incident electron and $\rho_{f}$ is the density of final states. The final state consists of the scattered electron, the emitted photon and the nucleus which is infinitely heavy.

$$
\begin{equation*}
\rho_{f}=\frac{1}{(2 \pi)^{6}} E_{2} p_{2} d \Omega_{2} \omega^{2} d \omega d \Omega_{\omega}, \tag{6.162}
\end{equation*}
$$

where $E_{2}, p_{2}$ denote the energy and momentum of the scattered electron, $\omega$, the frequency of the emitted photon which has a continuous spectrum. The solid angles $d \Omega_{2}$ and $d \Omega_{\omega}$ correspond to those of scattered electron and emitted photon.

The evaluation of the square of the matrix element is similar to that of Compton scattering and only the final result for the differential cross section is given below ${ }^{15}$ :

$$
\begin{align*}
d \sigma= & \frac{1}{2 \pi}\left(\frac{Z e^{2}}{q^{2}}\right)^{2} e^{2} \frac{d \omega}{\omega} \frac{p_{2}}{p_{1}} \sin \theta_{2} d \theta_{2} \sin \theta_{1} d \theta_{1} d \phi \\
& \times\left\{\frac{p_{2}^{2} \sin ^{2} \theta_{2}\left(4 E_{1}^{2}-q^{2}\right)}{\left(E_{2}-p_{2} \cos \theta_{2}\right)^{2}}+\frac{p_{1}^{2} \sin ^{2} \theta_{1}\left(4 E_{2}^{2}-q^{2}\right)}{\left(E_{1}-p_{1} \cos \theta_{1}\right)^{2}}\right. \\
& \left.-\frac{2 A p_{1} p_{2} \sin \theta_{1} \sin \theta_{2} \cos \phi-2 B \omega^{2}}{\left(E_{2}-p_{2} \cos \theta_{2}\right)\left(E_{1}-p_{1} \cos \theta_{1}\right)}\right\}, \tag{6.163}
\end{align*}
$$

with

$$
A=4 E_{1} E_{2}-q^{2}+2 \omega^{2}, \quad B=p_{2}^{2} \sin ^{2} \theta_{2}+p_{1}^{2} \sin ^{2} \theta_{1} .
$$

[^48]
### 6.6.1 In the soft photon limit

If the emitted photon energy is small compared to the rest mass of the electron (in the soft photon limit), then an approximate expression for the Bremsstrahlung cross section which has a simple interpretation can be obtained.

Writing the matrix element in terms of $\nless k$ instead of $q$, we have

$$
\begin{align*}
\mathcal{M}=-i\left(4 \pi e^{2}\right)^{1 / 2} & {\left[\bar{u}\left(\mathbf{p}_{2}\right) \notin(\mathbf{k}) \frac{1}{\not p_{2}+\nmid-m} \phi(\mathbf{q}) u\left(\mathbf{p}_{1}\right)\right.} \\
& \left.+\bar{u}\left(\mathbf{p}_{2}\right) \phi(\mathbf{q}) \frac{1}{\not p_{1}-\not b-m} \phi(\mathbf{k}) u\left(\mathbf{p}_{1}\right)\right], \tag{6.164}
\end{align*}
$$

Rationalizing the denominator of the matrix element, we get

$$
\begin{align*}
& \mathcal{M}=-i\left(4 \pi e^{2}\right)^{1 / 2} \frac{4 \pi z e^{2}}{q^{2}} \bar{u}\left(\mathbf{p}_{2}\right)\left[\notin(\mathbf{k}) \frac{\not p_{2}+\not \hbar+m}{\left(\not p_{2}+\not \vDash\right)^{2}-m^{2}} \gamma_{0}\right. \\
& \left.+\gamma_{0} \frac{\not p_{1}-\not \nmid k+m}{\left.\left(\not p_{1}-\not \not\right)^{2}\right)^{2}-m^{2}} \notin(\mathbf{k})\right] u\left(\mathbf{p}_{1}\right) \\
& =-i\left(4 \pi e^{2}\right)^{1 / 2} \frac{4 \pi z e^{2}}{q^{2}} \bar{u}\left(\mathbf{p}_{2}\right)\left[\notin(\mathbf{k}) \frac{\not p_{2}+\not \neq+m}{2 \mathbf{p}_{2} \cdot \mathbf{k}} \gamma_{0}\right. \\
& \left.+\gamma_{0} \frac{\not p_{1}-\not p+m}{\left(-2 \mathbf{p}_{1} \cdot \mathbf{k}\right)} \notin(\mathbf{k})\right] u\left(\mathbf{p}_{1}\right), \tag{6.165}
\end{align*}
$$

since

$$
\begin{align*}
& \left(\not p_{2}+\not \nmid\right)^{2}-m^{2}=\mathbf{p}_{2}^{2}+\mathbf{k}^{2}+2 \mathbf{p}_{2} \cdot \mathbf{k}-m^{2}=2 \mathbf{p}_{2} \cdot \mathbf{k}  \tag{6.166}\\
& \left.\left(\not p_{1}-\not \not\right)^{2}\right)^{2}-m^{2}=\mathbf{p}_{1}^{2}+\mathbf{k}^{2}-2 \mathbf{p}_{1} \cdot \mathbf{k}-m^{2}=-2 \mathbf{p}_{1} \cdot \mathbf{k} . \tag{6.167}
\end{align*}
$$

It is convenient to write the matrix element (6.165) in a simple form.

$$
\begin{equation*}
\mathcal{M}=-i\left(4 \pi e^{2}\right)^{1 / 2} \frac{4 \pi z e^{2}}{q^{2}}\left\{\frac{T_{a}}{2 \mathbf{p}_{2} \cdot \mathbf{k}}-\frac{T_{b}}{2 \mathbf{p}_{1} \cdot \mathbf{k}}\right\} \tag{6.168}
\end{equation*}
$$

with

$$
\begin{align*}
T_{a} & =\bar{u}\left(\mathbf{p}_{2}\right) \notin(\mathbf{k})\left(\not p_{2}+\not b+m\right) \gamma_{0} u\left(\mathbf{p}_{1}\right)  \tag{6.169}\\
T_{b} & =\bar{u}\left(\mathbf{p}_{2}\right) \gamma_{0}\left(\not p_{1}-\not b+m\right) \notin(\mathbf{k}) u\left(\mathbf{p}_{1}\right) \tag{6.170}
\end{align*}
$$

In the soft photon limit, the $\not \not k t$ term in $T_{a}$ and $T_{b}$ can be neglected. Then

$$
\begin{align*}
T_{a} & =\bar{u}\left(\mathbf{p}_{2}\right)\left\{\notin(\mathbf{k})\left(\not p_{2}+m\right) \gamma_{0}\right\} u\left(\mathbf{p}_{1}\right) \\
& =\bar{u}\left(\mathbf{p}_{2}\right)\left(-\not p_{2} \notin(\mathbf{k})+2 \mathbf{p}_{2} \cdot \boldsymbol{\varepsilon}+m \notin(\mathbf{k})\right) \gamma_{0} u\left(\mathbf{p}_{1}\right) \\
& =\bar{u}\left(\mathbf{p}_{2}\right) \gamma_{0} u\left(\mathbf{p}_{1}\right) 2 \mathbf{p}_{2} \cdot \boldsymbol{\varepsilon}, \quad \text { since } \bar{u}\left(\mathbf{p}_{2}\right) \not p_{2}=m \bar{u}\left(\mathbf{p}_{2}\right) . \tag{6.171}
\end{align*}
$$

In a similar way, we get

$$
\begin{align*}
T_{b} & =\bar{u}\left(\mathbf{p}_{2}\right)\left\{\gamma_{0}\left(\not p_{1}+m\right) \notin(\mathbf{k})\right\} u\left(\mathbf{p}_{1}\right) \\
& =\bar{u}\left(\mathbf{p}_{2}\right) \gamma_{0}\left(-\notin(\mathbf{k}) \not p_{1}+2 \mathbf{p}_{1} \cdot \boldsymbol{\varepsilon}+m \notin(\mathbf{k})\right) u\left(\mathbf{p}_{1}\right) \\
& =\bar{u}\left(\mathbf{p}_{2}\right) \gamma_{0} u\left(\mathbf{p}_{1}\right) 2 \mathbf{p}_{1} \cdot \boldsymbol{\varepsilon}, \quad \text { since } \not p_{1} u\left(\mathbf{p}_{1}\right)=m u\left(\mathbf{p}_{1}\right) . \tag{6.172}
\end{align*}
$$

Substituting these results in Eq. (6.168), we get

$$
\begin{equation*}
\mathcal{M}=-i\left(4 \pi e^{2}\right)^{1 / 2} \frac{4 \pi z e^{2}}{q^{2}}\left(\bar{u}\left(\mathbf{p}_{2}\right) \gamma_{0} u\left(\mathbf{p}_{1}\right)\right)\left(\frac{\mathbf{p}_{2} \cdot \boldsymbol{\varepsilon}(\mathbf{k})}{\mathbf{p}_{2} \cdot \mathbf{k}}-\frac{\mathbf{p}_{1} \cdot \boldsymbol{\varepsilon}(\mathbf{k})}{\mathbf{p}_{1} \cdot \mathbf{k}}\right)(, \in \tag{6.173}
\end{equation*}
$$

Squaring the matrix element and taking the sum over the final spin states and the average over the initial spin states of the electron, we get

$$
\begin{align*}
d \sigma= & \frac{2 \pi}{v_{1}} \frac{1}{2 E_{1} 2 E_{2} 2 \omega}|\overline{\mathcal{M}}|^{2} \frac{E_{2} p_{2} d \Omega_{2}}{(2 \pi)^{3}} \frac{\omega^{2} d \omega d \Omega_{\omega}}{(2 \pi)^{3}} \\
= & {\left[\frac{2 \pi}{v_{1}} \frac{1}{2 E_{1} 2 E_{2}}\left(\frac{4 \pi z e^{2}}{q^{2}}\right)^{2} \sum_{\text {spins }}\left|\bar{u}\left(\mathbf{p}_{2}\right) \gamma_{0} u\left(\mathbf{p}_{1}\right)\right|^{2} \frac{E_{2} p_{2} d \Omega_{2}}{(2 \pi)^{3}}\right] } \\
& \times\left[\frac{\left(4 \pi e^{2}\right) \omega^{2} d \Omega_{\omega}}{2 \omega(2 \pi)^{3}}\left(\frac{\mathbf{p}_{2} \cdot \varepsilon(\mathbf{k})}{\mathbf{p}_{2} \cdot \mathbf{k}}-\frac{\mathbf{p}_{1} \cdot \varepsilon(\mathbf{k})}{\mathbf{p}_{1} \cdot \mathbf{k}}\right)^{2}\right] \tag{6.174}
\end{align*}
$$

where the first square bracket denotes the elastic scattering cross section of electron in the Coulomb field making a transition from an initial state 1 to a final state 2 and the second square bracket gives the probability of emission of photon in frequency range $d \omega$ and solid angle $d \Omega_{\omega}$.

### 6.7 Electron-positron pair production

A single photon of energy greater than twice the electron mass cannot create an electron-positron pair since energy and momentum has to be conserved at each vertex and also for the whole process. Although two photons can create a pair, this process is extremely unlikely since the photon density is generally very low. However, a single photon can create a pair in the Coulomb field of a nucleus by sharing some momentum as shown in Fig. 6.10. It can happen in two ways: (a) The produced positron (virtual) can share momentum with the external field and emerge as a real positron or (b) the produced electron (virtual) can share momentum with the external field and emerge as real electron. These two modes are

(a)

(b)

Figure 6.10: Electron-positron pair production in the Coulomb field of a nucleus by a single photon. The four-momenta satisfy the relation $\mathbf{p}_{1}+\mathbf{q}+\mathbf{k}=\mathbf{p}_{2}$.
indistinguishable and so the two amplitudes should be added to compute the cross section for the process.

By an inspection of the two Feynman diagrams (a) and (b), one can write down the total matrix element for the process.

$$
\begin{align*}
& \mathcal{M}=-i\left(4 \pi e^{2}\right)^{1 / 2} {\left[u\left(\mathbf{p}_{2}\right) \notin(\mathbf{k}) \frac{1}{\not p_{1}+\not q+m} \phi(\mathbf{q}) v\left(\mathbf{p}_{1}\right)\right.} \\
&+u\left(\mathbf{p}_{2}\right) \phi(\mathbf{q}) \frac{1}{\not p_{1}+\not \not-m} \bar{m}  \tag{6.175}\\
&\left.\phi(\mathbf{k}) v\left(\mathbf{p}_{1}\right)\right],
\end{align*}
$$

with

$$
\phi(\mathbf{q})=\frac{4 \pi z e^{2}}{q^{2}} \delta\left(\mathbf{q}_{0}\right) \gamma_{0} \quad \text { and } \quad \notin=\sum_{\mu} \epsilon_{\mu} \gamma_{\mu}
$$

The matrix elements for this process are similar to those of Bremsstrahlung and can be evaluated in a similar way.

### 6.8 Muon pair production in electron-positron collision

Hitherto, we have considered only processes involving electrons, positrons and photons. The formalism can easily be extended to include production of new particles in $e^{-}-e^{+}$collisions.

Let us consider the process

$$
e^{-}+e^{+} \longrightarrow \mu^{-}+\mu^{+} .
$$

This process is possible if the energy of the colliding partners $e^{-}$and $e^{+}$in the c.m. frame is greater than the mass equivalent of $\mu^{-}$and $\mu^{+}$i.e. 212 MeV . At such energies, one can neglect the electron mass ( 0.511 MeV ) in the calculation.


Figure 6.11: Muon production in $e^{-}, e^{+}$collision. The Feynman diagram (a) represents the process $e^{-} e^{+} \longrightarrow \mu^{-} \mu^{+}$. Diagram (b) represents the kinematics for the process in c.m. frame. $E_{1}, \boldsymbol{p}_{1}$ and $E_{2}, \boldsymbol{p}_{2}$ denote the energy and momentum of the colliding partners $e^{-}$and $e^{+}$and $E_{3}, \boldsymbol{p}_{3}$ and $E_{4}, \boldsymbol{p}_{4}$ denote the energy and momentum of the produced particles $\mu^{-}$and $\mu^{+} . \theta$ denotes the angle at which the produced particles emerge in the c.m. system.

Using the Feynman rules, the matrix element for the process $e^{-} e^{+} \longrightarrow$ $\mu^{-} \mu^{+}$can be written down. For writing down the matrix elements, we follow the flow direction of arrow marks in Feynman diagram. We follow the same convention that we used earlier and write down the spinors $u\left(\mathbf{p}_{1}\right), \bar{v}\left(\mathbf{x}_{2}\right), \bar{u}\left(\mathbf{x}_{3}\right),\left(\bar{v}\left(\mathbf{x}_{4}\right)\right.$ as $u_{1}, \bar{v}_{2}, \bar{u}_{3}, \bar{v}_{4}$.

$$
\begin{equation*}
\mathcal{M}=i 4 \pi e^{2} \frac{\left[\left(\bar{v}_{2} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{3} \gamma_{\mu} v_{4}\right)\right]}{\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{2}} \tag{6.176}
\end{equation*}
$$

where $u_{1}$ and $u_{3}$ denote the spinors corresponding to particles ( $e^{-}$and $\mu^{-}$) and $v_{2}$ and $v_{4}$ denote the spinors corresponding to antiparticles ( $e^{+}$ and $\mu^{+}$). The four momenta of the colliding electron and positron are denoted by $\mathbf{p}_{1}$ and $\mathbf{p}_{2}$ and the four-momenta of the produced $\mu^{-}$and $\mu^{+}$are denoted by $\mathbf{p}_{3}$ and $\mathbf{p}_{4}$. The transition probability is the absolute square of the matrix element (6.176). Since the spins are not observed, a sum over the final spin states and an average over the initial spin states
have to be taken. This introduces a factor $\frac{1}{4}$.

$$
\begin{align*}
|\overline{\mathcal{M}}|^{2} & =\frac{1}{4} \mathcal{M} \mathcal{M}^{\dagger} \\
& =\frac{4 \pi^{2} e^{4}}{\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{4}} \sum_{\text {spins }}\left[\left(\bar{v}_{2} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{3} \gamma_{\mu} v_{4}\right)\right]\left[\left(\bar{v}_{2} \gamma_{\nu} u_{1}\right)\left(\bar{u}_{3} \gamma_{\nu} v_{4}\right)\right]^{\dagger} \\
& =\frac{4 \pi^{2} e^{4}}{\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{4}} \sum_{\text {spins }}\left[\left(\bar{v}_{2} \gamma_{\mu} u_{1}\right)\left(\bar{v}_{2} \gamma_{\nu} u_{1}\right)^{\dagger}\right]\left[\left(\bar{u}_{3} \gamma_{\mu} v_{4}\right)\left(\bar{u}_{3} \gamma_{\nu} v_{4}\right)^{\dagger}\right] . \\
& =\frac{4 \pi^{2} e^{4}}{\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{4}} \sum_{\text {spins }}\left[\left(\bar{v}_{2} \gamma_{\mu} u_{1}\right)\left(\bar{u}_{1} \gamma_{\nu} v_{2}\right)\right]\left[\left(\bar{u}_{3} \gamma_{\mu} v_{4}\right)\left(\bar{v}_{4} \gamma_{\nu} u_{3}\right)\right] .(6.1 \tag{6.177}
\end{align*}
$$

Using the trace techniques, we obtain

$$
\begin{align*}
\left.|\overline{\mathcal{M}}|^{2}=\frac{4 \pi^{2} e^{4}}{\left(\mathbf{p}_{1}+\right.}+\mathbf{p}_{2}\right)^{4} & \operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{1}+m_{e}\right) \gamma_{\nu}\left(\not p_{2}-m_{e}\right)\right] \\
& \times \operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{4}-m_{\mu}\right) \gamma_{\nu}\left(\not p_{3}+m_{\mu}\right)\right], \tag{6.178}
\end{align*}
$$

where $m_{e}$ and $m_{\mu}$ denote the mass of the electron and muon respectively. Let us now evaluate the traces, neglecting the terms involving $m_{e}$, since the energy of the colliding $e^{-}$and $e^{+}$should be greater than twice the muon mass ( $\sim 212 \mathrm{MeV}$ ). Since the trace of a product of odd number of $\gamma$ matrices vanishes, we get

$$
\begin{align*}
\operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{1}+m_{e}\right) \gamma_{\nu}\left(\not p_{2}-m_{e}\right)\right]= & \operatorname{Tr}\left[\gamma_{\mu} \not p_{1} \gamma_{\nu} \not p_{2}\right] \\
= & 4 p_{1 \mu} p_{2 \nu}+4 p_{2 \mu} p_{1 \nu}-4\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right) g_{\mu \nu} ;  \tag{6.179}\\
\operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{4}-m_{\mu}\right) \gamma_{\nu}\left(\not p_{3}+m_{\mu}\right)\right]= & \operatorname{Tr}\left[\gamma_{\mu} \not p_{4} \gamma_{\nu} p_{3}\right]-m_{\mu}^{2} \operatorname{Tr}\left[\gamma_{\mu} \gamma_{\nu}\right] \\
= & 4 p_{3 \mu} p_{4 \nu}+4 p_{4 \mu} p_{3 \nu}-4\left(\mathbf{p}_{3} \cdot \mathbf{p}_{4}\right) g_{\mu \nu} \\
& -4 m_{\mu}^{2} g_{\mu \nu} . \tag{6.180}
\end{align*}
$$

Inserting (6.179) and (6.180) in Eq. (6.178), we get

$$
\begin{align*}
|\overline{\mathcal{M}}|^{2}= & \frac{4 \pi^{2} e^{4}}{\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{4}}\left[4 p_{1 \mu} p_{2 \nu}+4 p_{2 \mu} p_{1 \nu}-4\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right) g_{\mu \nu}\right] \\
& \times\left[4 p_{3 \mu} p_{4 \nu}+4 p_{4 \mu} p_{3 \nu}-4\left(\mathbf{p}_{3} \cdot \mathbf{p}_{4}\right) g_{\mu \nu}-4 m_{\mu}^{2} g_{\mu \nu}\right] \\
= & \frac{128 \pi^{2} e^{4}}{\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{4}}\left[\left(\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right)\left(\mathbf{p}_{2} \cdot \mathbf{p}_{4}\right)+\left(\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right)\left(\mathbf{p}_{2} \cdot \mathbf{p}_{3}\right)+m_{\mu}^{2}\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)\right] . \\
= & \frac{128 \pi^{2} e^{4}}{\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{4}}\left[\left(\mathbf{p}_{1} \cdot \mathbf{p}_{3}\right)^{2}+\left(\mathbf{p}_{1} \cdot \mathbf{p}_{4}\right)^{2}+m_{\mu}^{2}\left(\mathbf{p}_{1} \cdot \mathbf{p}_{2}\right)\right] . \tag{6.181}
\end{align*}
$$

The last step in Eq. (6.181) is obtained using the energy-momentum conservation in the reaction

$$
\mathbf{p}_{1}+\mathbf{p}_{2}=\mathbf{p}_{3}+\mathbf{p}_{4},
$$

which can be equally expressed as

$$
\begin{aligned}
& \mathbf{p}_{1}-\mathbf{p}_{3}=\mathbf{p}_{4}-\mathbf{p}_{2} \longrightarrow \mathbf{p}_{1} \cdot \mathbf{p}_{3}=\mathbf{p}_{2} \cdot \mathbf{p}_{4} . \\
& \mathbf{p}_{1}-\mathbf{p}_{4}=\mathbf{p}_{3}-\mathbf{p}_{2} \longrightarrow \mathbf{p}_{1} \cdot \mathbf{p}_{4}=\mathbf{p}_{2} \cdot \mathbf{p}_{3} .
\end{aligned}
$$

Let us evaluate $|\overline{\mathcal{M}}|^{2}$ given by Eq. (6.181) in c.m. frame using the kinematics described in Fig. (b).

$$
\begin{align*}
\mathbf{p}_{1} \cdot \mathbf{p}_{3} & =E_{1} E_{3}-\boldsymbol{p}_{1} \cdot \boldsymbol{p}_{3}=E^{2}-E p \cos \theta  \tag{6.182}\\
\mathbf{p}_{1} \cdot \mathbf{p}_{4} & =E_{1} E_{4}-\boldsymbol{p}_{1} \cdot \boldsymbol{p}_{4}=E^{2}+E p \cos \theta  \tag{6.183}\\
\mathbf{p}_{1} \cdot \mathbf{p}_{2} & =E_{1} E_{2}-\boldsymbol{p}_{1} \cdot \boldsymbol{p}_{2}=2 E^{2} ;  \tag{6.184}\\
\left(\mathbf{p}_{1}+\mathbf{p}_{2}\right)^{2} & =\mathbf{p}_{1}^{2}+\mathbf{p}_{2}^{2}+2 \mathbf{p}_{1} \cdot \mathbf{p}_{2} \approx 4 E^{2} \tag{6.185}
\end{align*}
$$

since

$$
\begin{aligned}
& E_{1}=E_{2}=E_{3}=E_{4}=E ; \quad \boldsymbol{p}_{2}=-\boldsymbol{p}_{1} ; \boldsymbol{p}_{4}=-\boldsymbol{p}_{3} ; \\
& \left|\boldsymbol{p}_{1}\right|=\left|\boldsymbol{p}_{2}\right|=E ; \quad\left|\boldsymbol{p}_{3}\right|=\left|\boldsymbol{p}_{4}\right|=p=\left(E^{2}-m_{\mu}^{2}\right)^{1 / 2} \\
& \mathbf{p}_{1}^{2}=\mathbf{p}_{2}^{2}=m_{e}^{2}, \quad \text { negligible when compared to their energy } E .
\end{aligned}
$$

Substituting (6.182) - (6.185) into Eq. (6.181) and simplifying, we get

$$
\begin{align*}
|\overline{\mathcal{M}}|^{2} & =\frac{8 \pi^{2} e^{4}}{E^{4}}\left[\left(E^{2}-E p \cos \theta\right)^{2}+\left(E^{2}+E p \cos \theta\right)^{2}+2 E^{2} m_{\mu}^{2}\right] \\
& =16 \pi^{2} e^{4}\left[\left(1+\cos ^{2} \theta\right)+\frac{m_{\mu}^{2}}{E^{2}}\left(1-\cos ^{2} \theta\right)\right] \tag{6.186}
\end{align*}
$$

The differential cross section is given by Fermi's golden rule

$$
d \sigma=\frac{2 \pi}{v}(\Pi N)^{-1}|\overline{\mathcal{M}}|^{2} \rho_{f},
$$

with

$$
v=2 \frac{p_{1}}{E_{1}} \approx 2 ; \quad(\Pi N)^{-1}=\frac{1}{2 E_{1} 2 E_{2} 2 E_{3} 2 E_{4}}=\frac{1}{16 E^{4}} ; \quad \rho_{f}=\frac{E_{3} p_{3}}{2(2 \pi)^{3}} d \Omega_{3} .
$$

Hence

$$
\begin{equation*}
\frac{d \sigma}{d \Omega_{3}}=\frac{e^{4}}{16 E^{2}}\left(1-\frac{m_{\mu}^{2}}{E^{2}}\right)^{1 / 2}\left[1+\cos ^{2} \theta+\frac{m_{\mu}^{2}}{E^{2}}\left(1-\cos ^{2} \theta\right)\right] \tag{6.187}
\end{equation*}
$$

Integrating over $d \Omega_{3}=2 \pi \sin \theta d \theta$, we get

$$
\int\left(1+\cos ^{2} \theta\right) d \Omega=\frac{16}{3} \pi ; \quad \int\left(1-\cos ^{2} \theta\right) d \Omega=\frac{8}{3} \pi ;
$$

which, in turn, yields the total cross section.

$$
\begin{align*}
\sigma & =\frac{\pi e^{4}}{3 E^{2}}\left(1-\frac{m_{\mu}^{2}}{E^{2}}\right)^{1 / 2}\left(1+\frac{m_{\mu}^{2}}{2 E^{2}}\right) \\
& =\frac{\pi \alpha^{2}}{3 E^{2}}\left(1-\frac{m_{\mu}^{2}}{E^{2}}\right)^{1 / 2}\left(1+\frac{m_{\mu}^{2}}{2 E^{2}}\right) \tag{6.188}
\end{align*}
$$

where $\alpha=e^{2}$ denotes the fine structure constant.

### 6.8.1 Application to hadron production

We have already calculated the cross section for the elastic scattering of $e^{-} e^{+}$. If the c.m. energy is sufficiently high, the inelastic events begin to contribute and if the energy exceeds several GeV , then several hadrons can also be produced. Thus the $e^{-} e^{+}$collision at high energies serve as a convenient tool for searching for new particles.

In the GeV range, one can neglect the muon mass and the total cross section for $e^{-} e^{+} \rightarrow \mu^{-} \mu^{+}$takes a simple form

$$
\begin{equation*}
\sigma=\frac{\pi \alpha^{2}}{3 E^{2}}=\frac{4 \pi \alpha^{2}}{3 s^{2}}, \tag{6.189}
\end{equation*}
$$

where $s=4 E^{2}$ denotes the square of total energy in c.m. system.
Equation (6.189) has been used successfully to determine the number of colour degree of freedom for the quarks. Since the hadrons are made up of quarks, the hadronic final state in $e^{-} e^{+}$collision occurs through the events $e^{-} e^{+} \rightarrow q_{i} \bar{q}_{i}$, where $q_{i}, \bar{q}_{i}$ denote the quark-antiquark pair of flavour $i$. If $e_{i}$ is the fractional charge of the quark $q_{i}$, then the production cross section for $e^{-} e^{+} \rightarrow q_{i} \bar{q}_{i}$ is

$$
\begin{equation*}
\sigma=\frac{\pi \alpha^{2}}{3 E^{2}} N_{c} e_{i}^{2} \tag{6.190}
\end{equation*}
$$

where $e_{i}$ denotes the fractional charge of quark of flavour $i$ and $N_{c}$ denotes the number of colours. The ratio

$$
\begin{equation*}
R=\frac{\sigma\left(e^{-} e^{+} \rightarrow \text { hadrons }\right)}{\sigma\left(e^{-} e^{+} \rightarrow \mu^{-} \mu^{+}\right)}=N_{c} \sum_{i} e_{i}^{2}, \tag{6.191}
\end{equation*}
$$

Restricting the quark flavours to $u, d, s, c$ and choosing $N_{c}=3$, the ratio of the cross sections $R$ is given in the following Table 6.1.

Table 6.1: Ratio of cross sections $R$ calculated with increasing c.m. collision energies indicating opening of new channels. Quark flavours: up $(u)$, down $(d)$, strange $(s), \operatorname{charm}(c) . N_{c}=3$.

| Quark flavour (i) | charge $\left(e_{i}\right)$ | $\sum_{i} e_{i}^{2}$ | $R=N_{c} \sum_{i} e_{i}^{2}$ |
| :---: | :---: | :---: | :---: |
| $u, d$ | $\frac{2}{3},-\frac{1}{3}$ | $\left(\frac{2}{3}\right)^{2}+\left(-\frac{1}{3}\right)^{2}=\frac{5}{9}$ | $\frac{5}{3}$ |
| $u, d, s$ | $\frac{2}{3},-\frac{1}{3},-\frac{1}{3}$ | $\frac{2}{3}$ | 2 |
| $u, d, s, c$ | $\frac{2}{3},-\frac{1}{3},-\frac{1}{3}, \frac{2}{3}$ | $\frac{10}{9}$ | $\frac{10}{3}$ |

The ( $e^{-} e^{+}$) colliding beam experiments performed in the energy range $2.5 \mathrm{GeV}-40 \mathrm{GeV}$ confirm the above theoretical prediction for $R$ and this is taken as the sufficient evidence for the three colour degrees of freedom for quarks. The sudden increase in the ratio of $R$ at certain energies points to the opening of new channels in the production of hadrons.

## Review Questions

6.1 Draw the Feynman diagram for the Coulomb scattering of an electron by a nucleus of charge $Z e$ and deduce the Mott scattering formula. Under what condition, it will reduce to the Rutherford scattering formula.
6.2 Draw the lowest order Feynman diagrams for the Compton scattering and write down the Feynman amplitudes for them. Deduce the Klein-Nishina formula. Show, how it reduces to the Thomson scattering formula in the non-relativistic limit.
6.3 Explain how you get two Feynman diagrams in the study of electronelectron scattering by one-photon exchange. Write down the Feynman diagrams and discuss how you can obtain the cross section for the process.
6.4 Draw the Feynman diagrams for electron-positron scattering and explain their significance. Discuss how you can obtain the differential cross section for this process.
6.5 Discuss the process of electron-positron pair annihilation into two photons and explain how you can determine the cross section for this process.
6.6 Discuss the process of emission of photons known as bremsstrahlung during the Coulomb scattering of an electron. For emission of soft photons, show the cross section can be expressed as a product of two factors, one giving the probability of Coulomb scattering and the other giving the probability of photon emission in a frequency interval.
6.7 Consider the process $e^{-} e^{+} \longrightarrow \mu^{-} \mu^{+}$and show that the total cross section for this process at very high energies is given by a simple expression $\sigma=$ $\frac{4 \pi \alpha^{2}}{3 s}$, where $\alpha$ denotes the fine structure constant and $s$ denotes the total c.m. energy of colliding particles.

## Problems

6.1 Using the anticommutation relations of $\gamma$ matrices, prove the following formulae involving a product of $\gamma$ matrices of repeated index ${ }^{16}$ :

$$
\begin{array}{ll}
\gamma_{\mu} \gamma_{\mu}=4 ; & \gamma_{\mu} \gamma_{\nu} \gamma_{\mu}=-2 \gamma_{\nu} \\
\gamma_{\mu} \gamma_{\nu} \gamma_{\rho} \gamma_{\mu}=4 g_{\nu \rho} ; & \gamma_{\mu} \gamma_{\nu} \gamma_{\rho} \gamma_{\lambda} \gamma_{\mu}=-2 \gamma_{\lambda} \gamma_{\rho} \gamma_{\nu}
\end{array}
$$

6.2 In the discussion on Compton scattering, we have come across the following traces as given by Eqs. (6.54) - (6.57):

$$
\begin{aligned}
T_{\mathrm{aa}} & =\operatorname{Tr}\left[\mathcal{O}_{\mathrm{a}}\left(\not p_{1}+m\right) \tilde{\mathcal{O}}_{\mathrm{a}}\left(\not p_{2}+m\right)\right] \\
T_{\mathrm{ab}} & =\operatorname{Tr}\left[\mathcal{O}_{\mathrm{a}}\left(\not p_{1}+m\right) \tilde{\mathcal{O}}_{\mathrm{b}}\left(\not p_{2}+m\right)\right] ; \\
T_{\mathrm{ba}} & =\operatorname{Tr}\left[\mathcal{O}_{\mathrm{b}}\left(\not{ }_{1}+m\right) \tilde{\mathcal{O}}_{\mathrm{a}}\left(\not p_{2}+m\right)\right] ; \\
T_{\mathrm{bb}} & =\operatorname{Tr}\left[\mathcal{O}_{\mathrm{b}}\left(\not{ }_{1}+m\right) \tilde{\mathcal{O}}_{\mathrm{b}}\left(\not p_{2}+m\right)\right] ;
\end{aligned}
$$

with the operators, explicitly given by Eqs. (6.58) - (6.61).

$$
\begin{aligned}
\mathcal{O}_{\mathrm{a}} & =\frac{1}{2 m \omega_{1}} k_{2} k_{1} \not \not k_{1} \\
\mathcal{O}_{\mathrm{b}} & =\frac{1}{2 m \omega_{2}} k_{1} k_{2} \not k_{2} \\
\tilde{\mathcal{O}}_{\mathrm{a}} & =\frac{1}{2 m \omega_{1}} \gamma_{0}\left(k_{2} k_{1} \not k_{1}\right)^{\dagger} \gamma_{0}=\frac{1}{2 m \omega_{1}}\left(\not k_{1} \xi_{1} k_{2}\right) \\
\tilde{\mathcal{O}}_{\mathrm{b}} & =\frac{1}{2 m \omega_{2}} \gamma_{0}\left(k_{1} k_{2} \not k_{2}\right)^{\dagger} \gamma_{0}=\frac{1}{2 m \omega_{2}}\left(\not k_{2} k_{2} k_{1}\right)
\end{aligned}
$$

[^49]Evaluate them and obtain the results given by Eqs. (6.62) - (6.65).

$$
\begin{aligned}
T_{\mathrm{aa}} & =\frac{2}{m \omega_{1}}\left\{m \omega_{2}+2\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)^{2}\right\} \\
T_{\mathrm{bb}} & =\frac{2}{m \omega_{2}}\left\{m \omega_{1}-2\left(\mathbf{k}_{2} \cdot \varepsilon_{1}\right)^{2}\right\} . \\
T_{\mathrm{ab}} & =2\left\{2\left(\varepsilon_{1} \cdot \varepsilon_{2}\right)^{2}-1\right\}-\frac{2}{m \omega_{1}}\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)^{2}+\frac{2}{m \omega_{2}}\left(\mathbf{k}_{2} \cdot \varepsilon_{1}\right)^{2} \\
T_{\mathrm{ba}} & =T_{\mathrm{ab}}
\end{aligned}
$$

6.3 Given the Klein-Nishina formula (6.73), obtain the total cross section for Compton scattering by performing the angular integration.
6.4 Draw the Feynman diagrams for positron-positron scattering and show that the differential cross section for this process is the same as that of electron-electron scattering..

## Solutions to Problems

6.1

$$
\gamma_{\mu} \gamma_{\mu}=\gamma_{0}^{2}-\gamma_{x}^{2}-\gamma_{y}^{2}-\gamma_{z}^{2}=4
$$

Using the relation $\gamma_{\mu} \gamma_{\nu}+\gamma_{\nu} \gamma_{\mu}=2 g_{\mu \nu}$ repeatedly, we get ${ }^{17}$

$$
\begin{aligned}
\gamma_{\mu} \underbrace{\gamma_{\nu} \gamma_{\mu}} & =\gamma_{\mu}\left(-\gamma_{\mu} \gamma_{\nu}+2 g_{\mu \nu}\right)=-4 \gamma_{\nu}+2 \gamma_{\nu}=-2 \gamma_{\nu} \\
\gamma_{\mu} \gamma_{\nu} \underbrace{\gamma_{\rho} \gamma_{\mu}} & =\gamma_{\mu} \gamma_{\nu}\left(-\gamma_{\mu} \gamma_{\rho}+2 g_{\mu \rho}\right) \\
& =-\gamma_{\mu} \underbrace{\gamma_{\nu} \gamma_{\mu}}_{\nu} \gamma_{\rho}+2 g_{\mu \rho} \gamma_{\mu} \gamma_{\nu} \\
& =\gamma_{\mu} \gamma_{\mu} \gamma_{\nu} \gamma_{\rho}-2 g_{\mu \nu} \gamma_{\mu} \gamma_{\rho}+2 g_{\mu \rho} \gamma_{\mu} \gamma_{\nu} \\
& =4 \gamma_{\nu} \gamma_{\rho}-2 \gamma_{\nu} \gamma_{\rho}+2 \gamma_{\rho} \gamma_{\nu} \\
& =2 \gamma_{\nu} \gamma_{\rho}+2 \gamma_{\rho} \gamma_{\nu}=4 g_{\nu \rho} . \\
\gamma_{\mu} \gamma_{\nu} \gamma_{\rho} \underbrace{\gamma_{\lambda} \gamma_{\mu}} & =\gamma_{\mu} \gamma_{\nu} \gamma_{\rho}\left(-\gamma_{\mu} \gamma_{\lambda}+2 g_{\mu \lambda}\right) \\
& =-\gamma_{\mu} \gamma_{\nu} \underbrace{\gamma_{\rho} \gamma_{\mu}} \gamma_{\lambda}+2 g_{\mu \lambda} \gamma_{\mu} \gamma_{\nu} \gamma_{\rho} \\
& =\gamma_{\mu} \underbrace{\gamma_{\nu} \gamma_{\mu} \gamma_{\rho} \gamma_{\lambda}-2 g_{\rho \mu} \gamma_{\mu} \gamma_{\nu} \gamma_{\lambda}+2 g_{\mu \lambda} \gamma_{\mu} \gamma_{\nu} \gamma_{\rho}} \\
& =-\gamma_{\mu} \gamma_{\mu} \gamma_{\nu} \gamma_{\rho} \gamma_{\lambda}+2 g_{\mu \nu} \gamma_{\mu} \gamma_{\rho} \gamma_{\lambda}-2 g_{\rho \mu} \gamma_{\mu} \gamma_{\nu} \gamma_{\lambda}+2 g_{\mu \lambda} \gamma_{\mu} \gamma_{\nu} \gamma_{\rho} \\
& =-4 \gamma_{\nu} \gamma_{\rho} \gamma_{\lambda}+2 \gamma_{\nu} \gamma_{\rho} \gamma_{\lambda}-2 \gamma_{\rho} \gamma_{\nu} \gamma_{\lambda}+2 \gamma_{\lambda} \gamma_{\nu} \gamma_{\rho} \\
& =-2 \gamma_{\nu} \gamma_{\rho} \gamma_{\lambda}-2 \gamma_{\rho} \gamma_{\nu} \gamma_{\lambda}+2 \gamma_{\lambda} \gamma_{\nu} \gamma_{\rho} \\
& =-4 g_{\nu \rho} \gamma_{\lambda}+2 \gamma_{\lambda} \gamma_{\nu} \gamma_{\rho}=-2 \gamma_{\lambda} \gamma_{\rho} \gamma_{\nu} .
\end{aligned}
$$

The last step is obtained using the relation $\gamma_{\nu} \gamma_{\rho}=-\gamma_{\rho} \gamma_{\nu}+2 g_{\nu \rho}$.

[^50]
### 6.2 To derive $T_{\mathrm{aa}}$

The following general relations are used in the evaluation of the traces. ${ }^{18}$

$$
\begin{align*}
& \phi b=-b \phi \phi+2 \mathbf{a} \cdot \mathbf{b}, \quad\left(\mathbf{a} \cdot \mathbf{b}=a_{\mu} b_{\mu}=a_{0} b_{0}-\boldsymbol{a} \cdot \mathbf{b}\right) ;  \tag{6.192}\\
& \phi \phi=\mathbf{a} \cdot \mathbf{a} ;  \tag{6.193}\\
& \operatorname{Tr}(\phi b)=4 \mathbf{a} \cdot \mathbf{b} ;  \tag{6.194}\\
& \operatorname{Tr}(d b \not b d)=4(\mathbf{a} \cdot \mathbf{b})(\mathbf{c} \cdot \mathbf{d})-4(\mathbf{a} \cdot \mathbf{c})(\mathbf{b} \cdot \mathbf{d})+4(\mathbf{a} \cdot \mathbf{d})(\mathbf{b} \cdot \mathbf{c}) ; \tag{6.195}
\end{align*}
$$

where $\mathbf{a}$ and $\mathbf{b}$ are any four vectors. The trace of a product of an odd number of $\gamma$ matrices vanishes. The trace of a product of $\gamma$ matrices remains the same if the product is taken in the reverse order.
The trace of a product of matrices $A, B, C$ is invariant under cyclic permutations and it obeys the distribution law.

$$
\begin{align*}
& \operatorname{Tr}(A B C)=\operatorname{Tr}(B C A)=\operatorname{Tr}(C A B) ;  \tag{6.196}\\
& \operatorname{Tr}(A+B+C)=\operatorname{Tr} A+\operatorname{Tr} B+\operatorname{Tr} C . \tag{6.197}
\end{align*}
$$

In the present case,

$$
\begin{align*}
\not p_{1} \not p_{1}=\mathbf{p}_{1} \cdot \mathbf{p}_{1}=m^{2} ; & \not k_{2} \not p_{2}=\mathbf{p}_{2} \cdot \mathbf{p}_{2}=m^{2} .  \tag{6.198}\\
\not k_{1} \not k_{1}=1=\mathbf{k}_{1} \cdot \mathbf{k}_{1}=0 ; & \not k_{2} \not k_{2}=\mathbf{k}_{2} \cdot \mathbf{k}_{2}=0 .  \tag{6.199}\\
\not k_{1} k_{1}= & -k_{1} \not k_{1}+2 \mathbf{k}_{1} \cdot \varepsilon_{1}=-k_{1} k_{1} ;  \tag{6.200}\\
\not k_{2} k_{2}= & -k_{2} \not k_{2}+2 \mathbf{k}_{2} \cdot \varepsilon_{2}=-k_{2} \not k_{2} ; \tag{6.201}
\end{align*}
$$

since $\mathbf{k}_{1} \cdot \varepsilon_{1}=0$ and $\mathbf{k}_{2} \cdot \varepsilon_{2}=0$ because $\boldsymbol{\varepsilon}_{1}$ and $\boldsymbol{\varepsilon}_{2}$ have no time components but their space components are normal to the momentum vectors $\boldsymbol{k}_{1}$ and $\boldsymbol{k}_{2}$ respectively. Also

$$
\begin{equation*}
\xi_{1} \xi_{1}=\varepsilon_{1} \cdot \varepsilon_{1}=-1 ; \quad \quad k_{2} \xi_{2}=\varepsilon_{2} \cdot \varepsilon_{2}=-1 \tag{6.202}
\end{equation*}
$$

In addition, if we choose the coordinate system to coincide with the rest system of the initial electron, then

$$
\begin{align*}
& p_{1} \xi_{1}=-k_{1} p_{1}+2 \mathbf{p}_{1} \cdot \varepsilon_{1}=-k_{1} \not p_{1},  \tag{6.203}\\
& \not p_{1} \xi_{2}=-k_{2} p_{1}+2 \mathbf{p}_{1} \cdot \varepsilon_{2}=-k_{2} p_{1}, \tag{6.204}
\end{align*}
$$

since $\mathbf{p}_{1} \cdot \varepsilon_{1}=\mathbf{p}_{1} \cdot \varepsilon_{2}=0$ because $\mathbf{p}_{1}$ has only the time component whereas $\varepsilon_{1}$ and $\varepsilon_{2}$ have only the space components.

$$
\begin{equation*}
T_{\mathrm{aa}}=\operatorname{Tr}\left[\mathcal{O}_{\mathrm{a}}\left(p_{1}+m\right) \tilde{\mathcal{O}}_{\mathrm{a}}\left(p_{2}+m\right)\right], \tag{6.205}
\end{equation*}
$$

[^51]where
\[

$$
\begin{align*}
& \mathcal{O}_{\mathrm{a}}=\frac{1}{2 m \omega_{1}}\left(\epsilon_{2} \xi_{1} \not k_{1}\right) \text {, } \\
& \tilde{\mathcal{O}}_{\mathrm{a}}=\frac{1}{2 m \omega_{1}} \gamma_{0}\left(k_{2} k_{1} \not k_{1}\right)^{\dagger} \gamma_{0}=\frac{1}{2 m \omega_{1}}\left(\not k_{1} \xi_{1} \xi_{2}\right) \text {. } \\
& T_{\mathrm{aa}}=\frac{1}{4 m^{2} \omega_{1}^{2}} \operatorname{Tr}\left[k_{2} k_{1} \not k_{1}\left(p_{1}+m\right) \not k_{1} \xi_{1} \xi_{2}\left(p_{2}+m\right)\right] \\
& =\frac{1}{4 m^{2} \omega_{1}^{2}} \operatorname{Tr}[\xi_{2} \xi_{1} \not k_{1} p_{1} \not k_{1} \xi_{1} \xi_{2} p_{2}+m^{2} k_{2} \xi_{1} \underbrace{l k_{1} k_{1}} k_{1} k_{1} \xi_{2}] . \tag{6.206}
\end{align*}
$$
\]

In the trace, we have omitted the cross terms involving $m$ since they involve odd number of $\gamma$ matrices and the trace of a product of an odd number of $\gamma$ matrices vanishes.
The $m^{2}$ term in the trace of Eq. (6.206) vanishes since ${ }^{19}$

$$
\begin{equation*}
\not k_{1} \nmid k_{1}=\mathbf{k}_{1} \cdot \mathbf{k}_{1}=0 \tag{6.207}
\end{equation*}
$$

and we are left with only one term.

$$
\begin{align*}
T_{\mathrm{aa}} & =\frac{1}{4 m^{2} \omega_{1}^{2}} \operatorname{Tr}[k_{2} \xi_{1} \underbrace{\not k_{1} p_{1}} k_{1} k_{1} \xi_{2} p_{2}] \\
& =\frac{1}{4 m^{2} \omega_{1}^{2}} \operatorname{Tr}\left[k_{2} k_{1}\left\{-p_{1} \not k_{1}+2\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right)\right\} \not k_{1} \xi_{1} k_{2} p_{2}\right], \tag{6.208}
\end{align*}
$$

using the relation

$$
\begin{equation*}
\not k_{1} p_{1}=-\not p_{1} \not k_{1}+2\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right) . \tag{6.209}
\end{equation*}
$$

Applying the relation (6.207) once again to Eq. (6.208), the first term vanishes and we are left with only the second term.

$$
\begin{align*}
T_{\mathrm{aa}} & =\frac{1}{4 m^{2} \omega_{1}^{2}} 2\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right) \operatorname{Tr}[k_{2} \underbrace{k_{1} \not k_{1}} \xi_{1} \xi_{2} \not p_{2}] \\
& =\frac{1}{4 m^{2} \omega_{1}^{2}} 2\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right) \operatorname{Tr}\left[k_{2}\left\{-\not k_{1} k_{1}+2\left(\mathbf{k}_{1} \cdot \varepsilon_{1}\right)\right\} \xi_{1} k_{2} b_{2}\right] \\
& =\frac{1}{4 m^{2} \omega_{1}^{2}} 2\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right) \operatorname{Tr}[-\xi_{2} \not k_{1} \underbrace{\left.k_{1} \xi_{1} k_{2} k_{2}\right]} \\
& =\frac{1}{4 m^{2} \omega_{1}^{2}} 2\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right) \operatorname{Tr}\left[k_{2} \not k_{1} \xi_{2} \not p_{2}\right] . \tag{6.210}
\end{align*}
$$

[^52]Equation (6.210) is obtained by using the relations ${ }^{20}$

$$
\begin{equation*}
\mathbf{k}_{1} \cdot \varepsilon_{1}=0 \quad \text { and } \quad \not \epsilon_{1} \xi_{1}=\varepsilon_{1} \cdot \varepsilon_{1}=-1 \tag{6.211}
\end{equation*}
$$

We shall evaluate the trace that occurs in Eq. (6.210).

$$
\begin{align*}
\operatorname{Tr}\left(\epsilon_{2} \not \hbar_{1} k_{2} \not p_{2}\right)= & 4\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)\left(\mathbf{p}_{2} \cdot \varepsilon_{2}\right)-4\left(\varepsilon_{2} \cdot \varepsilon_{2}\right)\left(\mathbf{p}_{2} \cdot \mathbf{k}_{1}\right) \\
& +4\left(\mathbf{p}_{2} \cdot \varepsilon_{2}\right)\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right) \\
= & 8\left(\mathbf{p}_{2} \cdot \varepsilon_{2}\right)\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)+4\left(\mathbf{p}_{2} \cdot \mathbf{k}_{1}\right) \tag{6.212}
\end{align*}
$$

Due to energy-momentum conservation,

$$
\begin{equation*}
\mathbf{p}_{2}=\mathbf{p}_{1}+\mathbf{k}_{1}-\mathbf{k}_{2} \tag{6.213}
\end{equation*}
$$

Consequently,

$$
\begin{equation*}
\mathbf{p}_{2} \cdot \varepsilon_{2}=\mathbf{p}_{1} \cdot \varepsilon_{2}+\mathbf{k}_{1} \cdot \varepsilon_{2}-\mathbf{k}_{2} \cdot \varepsilon_{2}=\mathbf{k}_{1} \cdot \varepsilon_{2} \tag{6.214}
\end{equation*}
$$

Substituting (6.212) and (6.214) into Eq. (6.210), we get

$$
\begin{equation*}
T_{\mathrm{aa}}=\frac{2}{m^{2} \omega_{1}^{2}}\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right)\left\{\mathbf{p}_{2} \cdot \mathbf{k}_{1}+2\left(\mathbf{k}_{1} \cdot \boldsymbol{\varepsilon}_{2}\right)^{2}\right\} \tag{6.215}
\end{equation*}
$$

Now, we shall choose the coordinate system in which the initial electron is at rest. Then

$$
\begin{align*}
\mathbf{p}_{1} \cdot \mathbf{k}_{1} & =m \omega_{1}  \tag{6.216}\\
\mathbf{p}_{2} \cdot \mathbf{k}_{1} & =\left(\mathbf{p}_{1}+\mathbf{k}_{1}-\mathbf{k}_{2}\right) \cdot \mathbf{k}_{1} \\
& =m \omega_{1}-\mathbf{k}_{1} \cdot \mathbf{k}_{2}, \operatorname{since} \mathbf{k}_{1} \cdot \mathbf{k}_{1}=0 \\
& =m \omega_{1}-\omega_{1} \omega_{2}(1-\cos \theta) \tag{6.217}
\end{align*}
$$

Using the compton frequency shift formula (6.35),

$$
1-\cos \theta=\frac{m}{\omega_{2}}-\frac{m}{\omega_{1}}
$$

we find

$$
\begin{equation*}
\mathbf{p}_{2} \cdot \mathbf{k}_{1}=m \omega_{2} \tag{6.218}
\end{equation*}
$$

Substituting (6.216) and (6.218) into Eq. (6.215), we get

$$
T_{\mathrm{aa}}=\frac{2}{m \omega_{1}}\left\{m \omega_{2}+2\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)^{2}\right\}
$$

which is the same as Eq. (6.62).

[^53]
## To derive $T_{\mathrm{bb}}$

This can be done by using the same procedure outlined above. But one can straight away write down the result using the crossing symmetry.

$$
\omega_{1} \longrightarrow \omega_{2} ; \quad \mathbf{k}_{1} \longrightarrow-\mathbf{k}_{2} ; \quad \varepsilon_{2} \longrightarrow \varepsilon_{1}
$$

This will yield the expression (6.63) for $T_{\mathrm{bb}}$.

To derive $T_{\mathrm{ab}}$

$$
\begin{equation*}
T_{\mathrm{ab}}=\operatorname{Tr}\left[\mathcal{O}_{\mathrm{a}}\left(\phi_{1}+m\right) \tilde{\mathcal{O}}_{\mathrm{b}}\left(p_{2}+m\right)\right], \tag{6.219}
\end{equation*}
$$

where

$$
\begin{align*}
& \mathcal{O}_{\mathrm{a}}=\frac{1}{2 m \omega_{1}}\left(\xi_{2} \xi_{1} \not k_{1}\right) ; \\
& \tilde{\mathcal{O}}_{\mathrm{b}}=\frac{1}{2 m \omega_{2}} \gamma_{0}\left(\xi_{1} \xi_{2} k_{2}\right)^{\dagger} \gamma_{0}=\frac{1}{2 m \omega_{2}}\left(k_{2} \xi_{2} \xi_{1}\right) \text {. } \\
& T_{\mathrm{ab}}=\frac{1}{4 m^{2} \omega_{1} \omega_{2}} \operatorname{Tr}\left[\xi_{2} \xi_{1} \not k_{1}\left(p_{1}+m\right) \not k_{2} \xi_{2} \xi_{1}\left(p_{2}+m\right)\right] \\
& =\frac{1}{4 m^{2} \omega_{1} \omega_{2}} \operatorname{Tr}\left[\xi_{2} k_{1} \not k_{1} p_{1} k_{2} k_{2} \xi_{2} \xi_{1} p_{2}+m^{2} k_{2} \xi_{1} \nmid k_{1} \not k_{2} \xi_{2} \xi_{1}\right] \\
& =\frac{1}{4 m^{2} \omega_{1} \omega_{2}}[A+B] \text {, } \tag{6.220}
\end{align*}
$$

where

$$
\begin{align*}
& A=\operatorname{Tr}\left[k_{2} k_{1} \not k_{1} \not p_{1} \psi_{2} k_{2} k_{2} k_{1} \not p_{2}\right] ;  \tag{6.221}\\
& B=m^{2} \operatorname{Tr}\left[\xi_{2} k_{1} \not k_{1} \not k_{2} k_{2} \xi_{2} k_{1}\right] . \tag{6.222}
\end{align*}
$$

In the trace, we have omitted the cross terms involving $m$ since they consist of odd number of $\gamma$ matrices and the trace of a product of an odd number of $\gamma$ matrices vanishes.
Due to energy-momentum conservation,

$$
\mathbf{p}_{2}=\mathbf{p}_{1}+\mathbf{k}_{1}-\mathbf{k}_{2} \quad \text { and } \quad \not p_{2}=\not p_{1}+\not \not 1_{1}-\not \not k_{2} .
$$

Substituting this in Eq. (6.221),

$$
\begin{align*}
A= & \operatorname{Tr}\left[\xi_{2} k_{1} \not k_{1} \not p_{1} \not k_{2} \xi_{2} k_{1} \not p_{1}\right]+\operatorname{Tr}\left[k_{2} \xi_{1} \not k_{1} p_{1} \nmid k_{2} \xi_{2} \xi_{1} \not k_{1}\right] \\
& -\operatorname{Tr}\left[\xi_{2} k_{1} k_{1} k_{1} b_{1} k_{2} \xi_{2} \xi_{2} k_{1} k_{2}\right]  \tag{6.223}\\
= & A_{1}+A_{2}+A_{3} .
\end{align*}
$$

First let us take $A_{1}$ and commute $\not p_{1}$ successively with the quantity on the right until it pairs with the other $\not p_{1}$. Using the result $\not p_{1} \cdot \not p_{1}=\mathbf{p}_{1}^{2}=m^{2}$, we find that one of the terms cancels with $B$.

$$
\begin{align*}
& A_{1}=\operatorname{Tr}[k_{2} k_{1} \not \psi_{1} \underbrace{\not p_{1} \nmid k_{2}} k_{2} k_{1} \not p_{1}] \\
& =-\operatorname{Tr}[k_{2} \xi_{1} \not k_{1} \not k_{2} \underbrace{\not{ }_{1} k_{2}} \xi_{1} \not p_{1}]+2\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right) \operatorname{Tr}[k_{2} \underbrace{\xi_{1} \not k_{1}} k_{2} k_{1} \not p_{1}] \\
& =\operatorname{Tr}[\xi_{2} k_{1} \not \not k_{1} \not \not k_{2} \xi_{2} \underbrace{p_{1} \xi_{1}} \not p_{1}]-2\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right) \operatorname{Tr}[\xi_{2} \not k_{1} \underbrace{\xi_{1} k_{2}} k_{1} \not p_{1}] \\
& =-\operatorname{Tr}[\xi_{2} k_{1} \not k_{1} \not k_{2} \xi_{2} k_{1} \underbrace{\not p_{1} \not p_{1}}]+2\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right) \operatorname{Tr}[k_{2} \not k_{1} \epsilon_{2} \underbrace{\xi_{1} \xi_{1}} \not p_{1}] \\
& -4\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right)\left(\varepsilon_{1} \cdot \varepsilon_{2}\right) \operatorname{Tr}\left[\epsilon_{2} \not k_{1} \xi_{1} \not p_{1}\right] \\
& =-m^{2} \operatorname{Tr}\left[k_{2} k_{1} \not k_{1} \not k_{2} k_{2} k_{1}\right]-2\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right) \operatorname{Tr}\left[k_{2} \not k_{1} k_{2} \not p_{1}\right] \\
& -4\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right)\left(\varepsilon_{1} \cdot \varepsilon_{2}\right) \operatorname{Tr}\left[\epsilon_{2} \not k_{1} \xi_{1} \not p_{1}\right] \\
& =-B-2\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right) \operatorname{Tr}\left[k_{2} \not k_{1} k_{2} p_{1}\right]-4\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right)\left(\varepsilon_{1} \cdot \varepsilon_{2}\right) \operatorname{Tr}\left[k_{2} \not k_{1} k_{1} p_{1}\right] \tag{6.224}
\end{align*}
$$

Thus

$$
\begin{aligned}
A_{1}+B= & -8\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right)\left\{\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)\left(\mathbf{p}_{1} \cdot \boldsymbol{\varepsilon}_{2}\right)-\left(\varepsilon_{2} \cdot \varepsilon_{2}\right)\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right)\right. \\
& \left.+\left(\mathbf{p}_{1} \cdot \varepsilon_{2}\right)\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)\right\}-16\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right)\left(\varepsilon_{1} \cdot \varepsilon_{2}\right)\left[\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)\left(\mathbf{p}_{1} \cdot \varepsilon_{1}\right)\right. \\
& \left.-\left(\varepsilon_{1} \cdot \boldsymbol{\varepsilon}_{2}\right)\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right)+\left(\mathbf{p}_{1} \cdot \varepsilon_{2}\right)\left(\mathbf{k}_{1} \cdot \varepsilon_{1}\right)\right] \\
= & 8\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right)\left(\varepsilon_{2} \cdot \varepsilon_{2}\right)\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right)+16\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right)\left(\varepsilon_{1} \cdot \varepsilon_{2}\right)^{2}\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right) \\
= & -8\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right)\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right)+16\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right)\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right)\left(\varepsilon_{1} \cdot \varepsilon_{2}\right)^{2},
\end{aligned}
$$

since $\varepsilon_{2} \cdot \varepsilon_{2}=-1$.

$$
\begin{aligned}
& A_{2}=\operatorname{Tr}\left[k_{2} \xi_{1} \not k_{1} \not b_{1} \not k_{2} k_{2} \xi_{1} \not k_{1}\right] \\
& =\operatorname{Tr}[\not k_{1} \underbrace{\not \not k_{1} \xi_{2}} \not \xi_{1} \not k_{1} \not b_{1} \not k_{2} k_{2}] \\
& =-\operatorname{Tr}[\not \xi_{1} k_{2} \underbrace{\not k_{1} k_{1}} \not k_{1} \not p_{1} \not k_{2} \xi_{2}]+2\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right) \operatorname{Tr}[\underbrace{\not \varepsilon_{1} k_{1}} \not k_{1} \not p_{1} \not k_{2} k_{2}]
\end{aligned}
$$

$$
\begin{align*}
& -2\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right) \operatorname{Tr}\left[\not k_{1} \not p_{1} \not k_{2} k_{2}\right] \\
& =-2\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right) \operatorname{Tr}\left[\not k_{1} \not p_{1} \not \not k_{2} \epsilon_{2}\right] \\
& =-8\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)\left[\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right)\left(\mathbf{k}_{2} \cdot \varepsilon_{2}\right)-\left(\mathbf{k}_{1} \cdot \mathbf{k}_{2}\right)\left(\mathbf{p}_{1} \cdot \boldsymbol{\varepsilon}_{2}\right)\right. \\
& \left.+\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right)\right] \\
& =-8\left(\mathbf{k}_{1} \cdot \boldsymbol{\varepsilon}_{2}\right)^{2}\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right) \tag{6.226}
\end{align*}
$$

$$
\begin{align*}
& A_{3}=-\operatorname{Tr}[\xi_{2} \xi_{1} \not \xi_{1} p_{1} \underbrace{l k_{2} k_{2}} \xi_{1} \xi_{1} \not \xi_{2}] \\
& =\operatorname{Tr}\left[\xi_{2} \xi_{1} \not k_{1} \not p_{1} \xi_{2} \nmid k_{2} \xi_{1} \nmid k_{2}\right]-2\left(\mathbf{k}_{2} \cdot \varepsilon_{2}\right) \operatorname{Tr}\left[\xi_{2} \xi_{1} \not k_{1} k_{1} k_{1} k_{1} k_{2}\right] \\
& =\operatorname{Tr}[\xi_{2} \xi_{1} \not k_{1} p_{1} \xi_{2} \underbrace{l k_{2} \xi_{1}} \nmid k_{2}] \\
& =-\operatorname{Tr}[k_{2} \xi_{1} k_{1} p_{1} \xi_{2} k_{1} \underbrace{l k_{2} k_{2}}]+2\left(\mathbf{k}_{2} \cdot \varepsilon_{1}\right) \operatorname{Tr}\left[k_{2} \xi_{1} \not k_{1} p_{1} \xi_{2} \not k_{2}\right] \\
& =2\left(\mathbf{k}_{2} \cdot \varepsilon_{1}\right) \operatorname{Tr}[k_{2} \xi_{1} \not k_{1} \not b_{1} \underbrace{k_{2} \not k_{2}}] \\
& =-2\left(\mathbf{k}_{2} \cdot \varepsilon_{1}\right) \operatorname{Tr}\left[\xi_{2} k_{1} k_{1} k_{1} p_{1} \not k_{2} \xi_{2}\right] \\
& =2\left(\mathbf{k}_{2} \cdot \varepsilon_{1}\right) \operatorname{Tr}\left[\epsilon_{1} \not k_{1} p_{1} \not k_{2}\right] \\
& =8\left(\mathbf{k}_{2} \cdot \varepsilon_{1}\right)\left[\left(\mathbf{k}_{1} \cdot \boldsymbol{\varepsilon}_{1}\right)\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right)-\left(\mathbf{p}_{1} \cdot \varepsilon_{1}\right)\left(\mathbf{k}_{1} \cdot \mathbf{k}_{2}\right)\right. \\
& \left.+\left(\mathbf{k}_{2} \cdot \varepsilon_{1}\right)\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right)\right] \\
& =8\left(\mathbf{k}_{2} \cdot \boldsymbol{\varepsilon}_{1}\right)^{2}\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right) \text {. } \tag{6.227}
\end{align*}
$$

Substituting (6.224) - (6.227) into Eq. (6.220), we get

$$
\begin{aligned}
T_{\mathrm{ab}}= & \frac{2}{m^{2} \omega_{1} \omega_{2}}\left[\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right)\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right)\left\{2\left(\varepsilon_{1} \cdot \varepsilon_{2}\right)^{2}-1\right\}-\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)^{2}\left(\mathbf{p}_{1} \cdot \mathbf{k}_{2}\right)\right. \\
& \left.+\left(\mathbf{k}_{2} \cdot \varepsilon_{1}\right)^{2}\left(\mathbf{p}_{1} \cdot \mathbf{k}_{1}\right)\right] \\
= & \frac{2}{m^{2} \omega_{1} \omega_{2}}\left[m^{2} \omega_{1} \omega_{2}\left\{2\left(\varepsilon_{1} \cdot \varepsilon_{2}\right)^{2}-1\right\}-m \omega_{2}\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)^{2}\right. \\
& \left.+m \omega_{1}\left(\mathbf{k}_{2} \cdot \varepsilon_{1}\right)^{2}\right] \\
= & 2\left\{2\left(\varepsilon_{1} \cdot \varepsilon_{2}\right)^{2}-1\right\}-\frac{2}{m \omega_{1}}\left(\mathbf{k}_{1} \cdot \varepsilon_{2}\right)^{2}+\frac{2}{m \omega_{2}}\left(\mathbf{k}_{2} \cdot \varepsilon_{1}\right)^{2} .
\end{aligned}
$$

This is the same expression as given by Eq. (6.64)
6.3 Let us start with the Klein-Nishina formula.

$$
\frac{d \sigma}{d \Omega}=\frac{e^{4}}{2 m^{2}}\left[\left(\frac{\omega_{2}}{\omega_{1}}\right)^{3}+\frac{\omega_{2}}{\omega_{1}}-\left(\frac{\omega_{2}}{\omega_{1}}\right)^{2}\left(1-\cos ^{2} \theta\right)\right] .
$$

Since $\omega_{2}$ depends on $\theta$ as shown in Eqs. (6.74) and (6.75), we need to take into account its dependence on $\theta$ while integrating. We shall perform the integration of the terms one by one.

## First term

$$
\int\left(\frac{\omega_{2}}{\omega_{1}}\right)^{3} d \Omega=2 \pi \int_{0}^{\pi} \frac{\sin \theta d \theta}{(1+r-r \cos \theta)^{3}}=2 \pi \int_{-1}^{+1} \frac{d x}{(a+b x)^{3}},
$$

with

$$
a=1+r, \quad b=-r, \quad x=\cos \theta
$$

Using the standard integral,

$$
\int \frac{d x}{(a+b x)^{3}}=-\frac{1}{2 b(a+b x)^{2}}
$$

and substituting the limits, we obtain

$$
\begin{equation*}
\int\left(\frac{\omega_{2}}{\omega_{1}}\right)^{3} d \Omega=4 \pi \frac{1+r}{(1+2 r)^{2}} \tag{6.228}
\end{equation*}
$$

## Second term

$$
\int\left(\frac{\omega_{2}}{\omega_{1}}\right) d \Omega=2 \pi \int_{0}^{\pi} \frac{\sin \theta d \theta}{1+r-r \cos \theta}=2 \pi \int_{-1}^{+1} \frac{d x}{a+b x}
$$

Using the standard integral,

$$
\int \frac{d x}{a+b x}=\frac{1}{b} \ln (a+b x)
$$

and substituting the limits, we obtain

$$
\begin{equation*}
\int\left(\frac{\omega_{2}}{\omega_{1}}\right) d \Omega=\frac{2 \pi}{r} \ln (1+2 r) \tag{6.229}
\end{equation*}
$$

## Third term

$$
\int\left(\frac{\omega_{2}}{\omega_{1}}\right)^{2} d \Omega=2 \pi \int_{0}^{\pi} \frac{\sin \theta d \theta}{(1+r-r \cos \theta)^{2}}=2 \pi \int_{-1}^{+1} \frac{d x}{(a+b x)^{2}}
$$

Using the standard integral,

$$
\int \frac{d x}{(a+b x)^{2}}=-\frac{1}{b(a+b x)}
$$

and substituting the limits, we obtain

$$
\begin{equation*}
\int\left(\frac{\omega_{2}}{\omega_{1}}\right)^{2} d \Omega=\frac{4 \pi}{1+2 r} \tag{6.230}
\end{equation*}
$$

## Fourth term

$$
\int\left(\frac{\omega_{2}}{\omega_{1}}\right)^{2} \cos ^{2} \theta d \Omega=2 \pi \int_{0}^{\pi} \frac{\cos ^{2} \theta \sin \theta d \theta}{(1+r-r \cos \theta)^{2}}=2 \pi \int_{-1}^{+1} \frac{x^{2} d x}{(a+b x)^{2}}
$$

Using the standard integral,

$$
\int \frac{x^{2} d x}{(a+b x)^{2}}=\frac{1}{b^{3}}\left\{a+b x-2 a \ln (a+b x)-\frac{a^{2}}{a+b x}\right\}
$$

and substituting the limits, we obtain

$$
\begin{equation*}
\int\left(\frac{\omega_{2}}{\omega_{1}}\right)^{2} \cos ^{2} \theta d \Omega=4 \pi\left\{\frac{1}{r^{2}}-\frac{1+r}{r^{3}} \ln (1+2 r)+\frac{(1+r)^{2}}{r^{2}(1+2 r)}\right\} \tag{6.231}
\end{equation*}
$$

Collecting all these results, we get the total cross section given by Eq. (6.76).
4.4 Feynman diagrams representing the positron-positron scattering is given below. Due to indistinguishability of positrons in the final state, one has to consider both diagrams (a) and (b).


Figure 6.12: Positron-positron scattering.

The matrix element for the process can easily be written down in the momentum representation by using the Feynman rules. One has to follow the direction of arrows in writing down the matrix elements. For the sake of brevity, we denote the spinors $v\left(\mathbf{p}_{1}\right), v\left(\mathbf{p}_{2}\right) \cdots$ by $v_{1}, v_{2}, \cdots$.

$$
\begin{equation*}
\mathcal{M}=i 4 \pi e^{2}\left[\frac{\left(\bar{v}_{4} \gamma_{\mu} v_{2}\right)\left(\bar{v}_{3} \gamma_{\mu} v_{1}\right)}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{2}}-\frac{\left(\bar{v}_{3} \gamma_{\mu} v_{2}\right)\left(\bar{v}_{4} \gamma_{\mu} v_{1}\right)}{\left(\mathbf{p}_{1}-\mathbf{p}_{4}\right)^{2}}\right], \tag{6.232}
\end{equation*}
$$

where summation over $\mu$ is implied. The first term in Eq. (6.232) corresponds to the diagram (a) and the second term corresponds to the exchange diagram (b). The negative sign between the two terms arises from the antisymmetry of the two-positron wave function. To find the transition probability, we have to find the absolute square of the matrix element (6.232). If the positron spins are not observed, then a sum over the final spin states and an average over the initial spin states have to be taken.

This introduces a factor $1 / 4$.

$$
\begin{aligned}
|\overline{\mathcal{M}}|^{2}= & \frac{1}{4} \mathcal{M} \mathcal{M}^{\dagger} \\
= & 4 \pi^{2} e^{4} \sum_{\mathrm{spins}}\left[\frac{\left(\bar{v}_{4} \gamma_{\mu} v_{2}\right)\left(\bar{v}_{3} \gamma_{\mu} v_{1}\right)}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{2}}-\frac{\left(\bar{v}_{3} \gamma_{\mu} v_{2}\right)\left(\bar{v}_{4} \gamma_{\mu} v_{1}\right)}{\left(\mathbf{p}_{1}-\mathbf{p}_{4}\right)^{2}}\right] \\
& \times\left[\frac{\left(\bar{v}_{4} \gamma_{\nu} v_{2}\right)\left(\bar{v}_{3} \gamma_{\nu} v_{1}\right)}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{2}}-\frac{\left(\bar{v}_{3} \gamma_{\nu} v_{2}\right)\left(\bar{v}_{4} \gamma_{\nu} v_{1}\right)}{\left(\mathbf{p}_{1}-\mathbf{p}_{4}\right)^{2}}\right]^{\dagger} \\
= & 4 \pi^{2} e^{4}\left[\frac{T_{\mathrm{aa}}}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{4}}+\frac{T_{\mathrm{bb}}}{\left(\mathbf{p}_{1}-\mathbf{p}_{4}\right)^{4}}-\frac{T_{\mathrm{ab}}+T_{\mathrm{ba}}}{\left(\mathbf{p}_{1}-\mathbf{p}_{3}\right)^{2}\left(\mathbf{p}_{1}-\mathbf{p}_{4}\right)^{2}}\right]
\end{aligned}
$$

where

$$
\begin{aligned}
T_{\mathrm{aa}} & =\sum_{\text {spins }}\left[\left(\bar{v}_{3} \gamma_{\mu} v_{1}\right)\left(\bar{v}_{3} \gamma_{\nu} v_{1}\right)^{\dagger}\left(\bar{v}_{4} \gamma_{\mu} v_{2}\right)\left(\bar{v}_{4} \gamma_{\nu} v_{2}\right)^{\dagger}\right] \\
& =\operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{1}-m\right) \gamma_{\nu}\left(\not p_{3}-m\right)\right] \operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{2}-m\right) \gamma_{\nu}\left(\not p_{4}-m\right)\right] \\
T_{\mathrm{bb}} & =\sum_{\text {spins }}\left[\left(\bar{v}_{4} \gamma_{\mu} v_{1}\right)\left(\bar{v}_{4} \gamma_{\nu} v_{1}\right)^{\dagger}\left(\bar{v}_{3} \gamma_{\mu} v_{2}\right)\left(\bar{v}_{3} \gamma_{\nu} v_{2}\right)^{\dagger}\right] \\
& =\operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{1}-m\right) \gamma_{\nu}\left(\not p_{4}-m\right)\right] \operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{2}-m\right) \gamma_{\nu}\left(\not p_{3}-m\right)\right], \\
T_{\mathrm{ab}} & =\sum_{\text {spins }}\left[\left(\bar{v}_{3} \gamma_{\mu} v_{1}\right)\left(\bar{v}_{4} \gamma_{\nu} v_{1}\right)^{\dagger}\left(\bar{v}_{4} \gamma_{\mu} v_{2}\right)\left(\bar{v}_{3} \gamma_{\nu} v_{2}\right)^{\dagger}\right] \\
& =\operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{1}-m\right) \gamma_{\nu}\left(\not p_{4}-m\right) \gamma_{\mu}\left(\not p_{2}-m\right) \gamma_{\nu}\left(\not p_{3}-m\right)\right] \\
T_{\mathrm{ba}} & =\sum_{\text {spins }}\left[\left(\bar{v}_{4} \gamma_{\mu} v_{1}\right)\left(\bar{v}_{3} \gamma_{\nu} v_{1}\right)^{\dagger}\left(\bar{v}_{3} \gamma_{\mu} v_{2}\right)\left(\bar{v}_{4} \gamma_{\nu} v_{2}\right)^{\dagger}\right] \\
& =\operatorname{Tr}\left[\gamma_{\mu}\left(\not p_{1}-m\right) \gamma_{\nu}\left(\not p_{3}-m\right) \gamma_{\mu}\left(\not p_{2}-m\right) \gamma_{\nu}\left(\not p_{4}-m\right)\right] .
\end{aligned}
$$

The expressions for $T_{\mathrm{aa}}, T_{\mathrm{bb}}, T_{\mathrm{ab}}, T_{\mathrm{ba}}$ given above have been obtained using the projection operators for positrons

$$
\sum_{\text {spins }} v \bar{v}=\not p-m
$$

and they are identical with the expressions (6.83) - (6.86) obtained earlier for electron-electron scattering except for the change in sign for $m$. Since the traces of a product of odd number of gamma matrices vanish, one has to consider the terms involving even number of gamma matrices only. That means one has to consider only products of even number of $m$ and so the results will be independent of the sign of $m$.
Hence, it follows that the positron-positron scattering cross section will be identical with electron-electron scattering cross section.

## Chapter 7

## Radiative Corrections

In the last chapter, we have studied the various processes in QED in the lowest order of perturbation theory. The higher order terms are expected to yield small corrections of the order of the fine structure constant $\alpha$. But on doing such a calculation, one encounters divergent integrals.

In the lowest order, the momenta of the particles in the intermediate stage are well defined because of the conservation of momenta at every vertex. But in higher order, the momenta of the intermediate particles can vary over a wide range and this causes divergences. Special techniques are to be devised to overcome these difficulties. These techniques have come to be known as Regularization and Renormalization. They are based on the arguments that what one considers in perturbation theory is the bare electron with mass $m_{0}$ and charge $-e_{0}$ and what one observes in experiments is the physical electron with mass $m$ and charge $-e$.

The higher order corrections are sometimes called Radiative corrections and they can be broadly classified into three types, electron self energy, photon self energy or vacuum polarization and vertex corrections. Any higher order correction to a given process will involve all these types and so it is sufficient to study them separately and devise methods to overcome the divergences.

### 7.1 Electron Self Energy

The Feynman diagram illustrating the mechanism by which an electron acquires self-energy is given in Fig. 7.1. The incident electron of mo-
mentum $\boldsymbol{p}$ emits a virtual photon of momentum $\boldsymbol{k}$ which is recaptured by the emerging electron of the same momentum $\boldsymbol{p}$. The virtual photon momentum $\boldsymbol{k}$ can assume any value from 0 to $\infty$ and hence an integration has to be performed over the momenta and this results in the divergence problem.


Figure 7.1: Self energy of the electron
The matrix element ${ }^{1}$ for this process in momentum representation can be written down ${ }^{2}$ using the Feynman rules.

$$
\begin{equation*}
\mathcal{M}=-i 4 \pi e^{2} \int \bar{u}(\boldsymbol{p}) \gamma_{\mu} \frac{1}{\not p-\not \vDash-m} \gamma_{\mu} u(\boldsymbol{p}) \frac{d^{4} k}{(2 \pi)^{4}} \frac{1}{\mathbf{k}^{2}} . \tag{7.1}
\end{equation*}
$$

Since the initial and final momenta of the electron are identical, the matrix element can be interpreted as the change in energy of the electron due to self-interaction. Introducing the normalization factor $\frac{1}{2 E}$ for the electron spinors, we obtain the following expression for the change in energy due to self interaction.

$$
\begin{equation*}
\Delta E=-i \frac{4 \pi e^{2}}{2 E} \int \bar{u}(\boldsymbol{p}) \gamma_{\mu} \frac{1}{\not p-\not p-m} \gamma_{\mu} u(\boldsymbol{p}) \frac{d^{4} k}{(2 \pi)^{4}} \frac{1}{\mathbf{k}^{2}} . \tag{7.2}
\end{equation*}
$$

Since the momentum of the electron remains the same after the perturbation, the change in energy can be interpreted as the change in mass of the electron.

$$
E^{2}=p^{2}+m^{2} ; \quad 2 E d E=2 m d m
$$

[^54]Thus, we obtain an expression for the change in mass of the electron due to self-interaction.

$$
\begin{equation*}
\Delta m=\frac{4 \pi e^{2}}{2 m i} \int \bar{u}\left(\gamma_{\mu} \frac{1}{\not p-\not p-m} \gamma_{\mu}\right) u \frac{d^{4} k}{(2 \pi)^{4}} \frac{1}{\mathbf{k}^{2}} . \tag{7.3}
\end{equation*}
$$

The quantity within the bracket can be simplified to yield

$$
\begin{equation*}
\gamma_{\mu} \frac{1}{\not p-\not p-m} \gamma_{\mu}=\frac{\gamma_{\mu}(\not p-\not p+m) \gamma_{\mu}}{(\not p-\not p)^{2}-m^{2}}=\frac{2 \not p+2 m}{\mathbf{k}^{2}-2 \mathbf{p} \cdot \mathbf{k}}, \tag{7.4}
\end{equation*}
$$

using the following relations.

$$
\begin{aligned}
& \gamma_{\mu} \not p \gamma_{\mu}=-2 \not p \rightarrow-2 m, \quad \text { since } \not p u=m u ; \\
& \gamma_{\mu} \not \not k \gamma_{\mu}=-2 \not b ; \quad \quad \gamma_{\mu} \gamma_{\mu}=4 ; \quad \not{ }^{2}=\mathbf{p}^{2}=m^{2} .
\end{aligned}
$$

Equation (7.3) now becomes

$$
\begin{equation*}
\Delta m=\frac{4 \pi e^{2}}{2 m i} \int \frac{\bar{u}(2 m+2 \not k) u}{\mathbf{k}^{2}-2 \mathbf{p} \cdot \mathbf{k}} \frac{d^{4} k}{(2 \pi)^{4}} \frac{1}{\mathbf{k}^{2}} . \tag{7.5}
\end{equation*}
$$

This integral is divergent and it has posed a major challenge in the study of QED for over 20 years. We shall closely follow Feynman's method of overcoming this problem. Let us modify the propagation kernel of the photon by multiplying it by a factor $c\left(\mathbf{k}^{2}\right)$.

$$
\begin{equation*}
\frac{1}{\mathbf{k}^{2}} \rightarrow \frac{1}{\mathbf{k}^{2}} c\left(\mathbf{k}^{2}\right)=\frac{1}{\mathbf{k}^{2}}\left(\frac{-\lambda^{2}}{\mathbf{k}^{2}-\lambda^{2}}\right)=\frac{1}{\mathbf{k}^{2}}-\frac{1}{\mathbf{k}^{2}-\lambda^{2}} \tag{7.6}
\end{equation*}
$$

The second term corresponds to the propagation of photon of mass $\lambda$.
Introducing this kernel into Eq. (7.5), we get

$$
\begin{equation*}
\Delta m=\frac{4 \pi e^{2}}{2 m i} \int \frac{\bar{u}(2 m+2 \nmid) u}{\mathbf{k}^{2}-2 \mathbf{p} \cdot \mathbf{k}} \frac{d^{4} k}{(2 \pi)^{4}} \frac{1}{\mathbf{k}^{2}}\left(\frac{-\lambda^{2}}{\mathbf{k}^{2}-\lambda^{2}}\right) . . \tag{7.7}
\end{equation*}
$$

A more convenient representation for the modified kernel is

$$
\begin{equation*}
-\int_{0}^{\lambda^{2}} \frac{d L}{\left(\mathbf{k}^{2}-L\right)^{2}}=-\left.\frac{1}{\mathbf{k}^{2}-L}\right|_{L=0} ^{L=\lambda^{2}}=-\frac{1}{\mathbf{k}^{2}-\lambda^{2}}+\frac{1}{\mathbf{k}^{2}} \tag{7.8}
\end{equation*}
$$

Using (7.8), we can rewrite Eq. (7.7) in the following form.

$$
\begin{equation*}
\Delta m=-\frac{4 \pi e^{2}}{2 m i} \int_{0}^{\lambda^{2}} d L \int \frac{\bar{u}(2 m+2 \not k) u}{\left(\mathbf{k}^{2}-2 \mathbf{p} \cdot \mathbf{k}\right)\left(\mathbf{k}^{2}-L\right)^{2}} \frac{d^{4} k}{(2 \pi)^{4}} . \tag{7.9}
\end{equation*}
$$

The denominator in Eq. (7.9) can be expressed as

$$
\begin{equation*}
\frac{1}{\left(\mathbf{k}^{2}-2 \mathbf{p} \cdot \mathbf{k}\right)\left(\mathbf{k}^{2}-L\right)^{2}}=\int_{0}^{1} \frac{2(1-x) d x}{\left[\mathbf{k}^{2}-2 x \mathbf{p} \cdot \mathbf{k}-L(1-x)\right]^{3}}, \tag{7.10}
\end{equation*}
$$

using the remarkable Feynman's formula (vide Solved Problem 5.2)

$$
\begin{equation*}
\frac{1}{a b^{2}}=\int_{0}^{1} \frac{2(1-x) d x}{[a x+b(1-x)]^{2}}, \tag{7.11}
\end{equation*}
$$

with the substitution

$$
a=\mathbf{k}^{2}-2 \mathbf{p} \cdot \mathbf{k}, \quad b=\mathbf{k}^{2}-L
$$

Substituting (7.10) into Eq. (7.9), we get

$$
\begin{equation*}
\Delta m=-\frac{4 \pi e^{2}}{2 m i} \int_{0}^{\lambda^{2}} d L \int_{0}^{1} d x \int \frac{d^{4} k}{(2 \pi)^{4}} \frac{2(1-x) \bar{u}(2 m+2 \not x) u}{\left[\mathbf{k}^{2}-2 x \mathbf{p} \cdot \mathbf{k}-L(1-x)\right]^{3}} . \tag{7.12}
\end{equation*}
$$

There are, in total, three integrations and let us perform them one by one.

## Integration over $d^{4} k$

The $d^{4} k$ integration can be performed using the following formula (vide Solved Problems (7.3) and (7.4)),

$$
\begin{equation*}
\int \frac{\left(1, k_{\sigma}\right)}{\left(\mathbf{k}^{2}-2 \mathbf{p} \cdot \mathbf{k}-\Delta\right)^{3}} \frac{d^{4} k}{(2 \pi)^{4}}=\frac{\left(1, p_{\sigma}\right)}{32 \pi^{2} i\left(\mathbf{p}^{2}+\Delta\right)} \tag{7.13}
\end{equation*}
$$

which is given in a compact form. The notation $\left(1, k_{\sigma}\right)$ on the L.H.S. means either 1 or $k_{\sigma}$ and it correspondingly yields either 1 or $p_{\sigma}$ on the R.H.S. represented by the notation ( $1, p_{\sigma}$ ).

By making the substitution

$$
\mathbf{p} \rightarrow x \mathbf{p} \quad \text { and } \quad \Delta \rightarrow L(1-x)
$$

we get

$$
\begin{align*}
\int \frac{\left(1, k_{\sigma}\right)}{\left[\mathbf{k}^{2}-2 x \mathbf{p} \cdot \mathbf{k}-L(1-x)\right]^{3}} \frac{d^{4} k}{(2 \pi)^{4}} & =\frac{\left(1, x p_{\sigma}\right)}{32 \pi^{2} i\left[x^{2} \mathbf{p}^{2}+L(1-x)\right]} \\
& =\frac{\left(1, x p_{\sigma}\right)}{32 \pi^{2} i\left[m^{2} x^{2}+L(1-x)\right]} \tag{7.14}
\end{align*}
$$

since $\mathbf{p}^{2}=m^{2}$.

## Integration over $d L$

Taking note of the factors that depend upon $L$, we get

$$
\begin{align*}
\int_{0}^{\lambda^{2}} \frac{d L}{m^{2} x^{2}+L(1-x)} & =\left.\frac{\ln \left[m^{2} x^{2}+L(1-x)\right]}{1-x}\right|_{L=0} ^{L=\lambda^{2}} \\
& =\frac{1}{1-x} \ln \left[\frac{m^{2} x^{2}+\lambda^{2}(1-x)}{m^{2} x^{2}}\right] \tag{7.15}
\end{align*}
$$

When $\lambda^{2} \gg m^{2}, m^{2} x^{2}$ occurring in the numerator can be neglected. Of course, when $x \approx .1$, $(1-x) \lambda^{2}$ is not much larger than $m^{2} x^{2}$, but the interval over which this is true is so small, for $\lambda^{2} \gg m^{2}$, that the error is small when the $x$-integration is performed. Under this approximation,

$$
\begin{equation*}
\int_{0}^{\lambda^{2}} \frac{d L}{m^{2} x^{2}+L(1-x)} \approx \frac{1}{1-x} \ln \left[\frac{\lambda^{2}(1-x)}{m^{2} x^{2}}\right] \tag{7.16}
\end{equation*}
$$

Substituting (7.14) and (7.16) into Eq. (7.12), we get

$$
\begin{align*}
\Delta m & =-\frac{4 \pi e^{2}}{2 m i} \frac{1}{16 \pi^{2} i} \int_{0}^{1} d x\{\bar{u}(2 m+2 x \not p) u\} \ln \left(\frac{\lambda^{2}(1-x)}{m^{2} x^{2}}\right) \\
& =\frac{e^{2}}{8 \pi m} \int_{)}^{1} d x\left\{\bar{u}(2 m+2 x \not p) u \ln \left(\frac{\lambda^{2}(1-x)}{m^{2} x^{2}}\right)\right. \\
& =\frac{e^{2} m}{2 \pi} \int_{0}^{1} d x(1+x) \ln \left(\frac{\lambda^{2}(1-x)}{m^{2} x^{2}}\right) . \tag{7.17}
\end{align*}
$$

To obtain the last step, we have used the relations $\not p u=m u$ and $\bar{u} u=2 m$.

## Integration over $d x$

Expanding the $\ln$ term

$$
\ln \left(\frac{\lambda^{2}(1-x)}{m^{2} x^{2}}\right)=\ln \frac{\lambda^{2}}{m^{2}}+\ln \frac{1-x}{x^{2}}=2 \ln \frac{\lambda}{m}+\ln \frac{1-x}{x^{2}},
$$

and using the definite integrals (vide Solved Problem 5.5)

$$
\begin{equation*}
\int_{0}^{1} \ln \left(\frac{1-x}{x^{2}}\right) d x=1 ; \quad \int_{0}^{1} x \ln \left(\frac{1-x}{x^{2}}\right) d x=-1 / 4 \tag{7.18}
\end{equation*}
$$

we get a simple expression for the change in mass due to self interaction.

$$
\begin{equation*}
\frac{\Delta m}{m}=\frac{e^{2}}{2 \pi}\left[3 \ln \frac{\lambda}{m}+\frac{3}{4}\right] . \tag{7.19}
\end{equation*}
$$

The factor $\left(e^{2} / 2 \pi\right)$ is of the order of $10^{-3}$. Even if $\lambda$ is many times $m$, the fractional change in mass will not be very large. However, the shift in mass depends upon the parameter $\lambda$ and hence cannot be determined theoretically. The theoretical and the experimental mass ${ }^{3}$ are related by

$$
m_{\mathrm{exp}}=m_{\mathrm{th}}+\delta m
$$

All our experimental measurements are made with experimental mass with self-action included and the theoretical mass (mass without selfaction) cannot be determined. The replacement of theoretical mass $m_{\text {th }}$ by the experimental mass $m_{\exp }$ is what is known as mass renormalization and $m_{\text {exp }}$ is often referred to as the renormalized mass.

For a free particle, a theory using $m_{\text {th }}$ plus self-action is equivalent to a theory using $m_{\exp }+$ self-action minus $\Delta m$. For a bound-state electron, a theory using $m_{\text {exp }}$ will slightly differ from this principle of equivalence and this leads to the Lamb shift in the hydrogen atom.

### 7.2 Higher order corrections

For the sake of illustration, let us consider the scattering of an electron by a static electric potential such as the Coulomb scattering by a nucleus of charge $z e$. Figure 7.2 represents the lowest order Feynman diagram.


Figure 7.2: Lowest order Feynman diagram representing the scattering of an electron by a static external potential (for example, Coulomb potential by a nucleus of charge $Z e$ ). $\boldsymbol{p}_{i}$ and $\boldsymbol{p}_{f}$ denote the momenta of the incoming and scattered electron and $\boldsymbol{q}$, the three-momentum transfer to the electron.

The Feynman amplitude corresponding to the lowest order diagram 5.2 is given by

$$
\begin{equation*}
\mathcal{M}=-i e \bar{u}\left(\boldsymbol{p}_{f}\right) \phi(\mathbf{q}) u\left(\boldsymbol{p}_{i}\right), \quad \boldsymbol{q}=\boldsymbol{p}_{f}-\boldsymbol{p}_{i} . \tag{7.20}
\end{equation*}
$$

[^55]The next higher order Feynman diagram for this process involves two additional vertices representing the emission and reabsorption of virtual photon as shown in fig. 7.3.


Figure 7.3: Radiative corrections to the scattering of an electron by a static external potential. $\boldsymbol{p}_{i}$ and $\boldsymbol{p}_{f}$ denote the momenta of the incoming and scattered electron and $\boldsymbol{q}$, the three-momentum transfer to the electron.

Restricting oneself to the Feynman diagrams which contain two extra vertices, involving only one virtual photon, we obtain the radiative corrections to the process of scattering of electron by the static electric field. As shown in Fig. 7.3, there are four possibilities that correspond to second order corrections to the basic process represented by Fig. 7.1. Each of them can be considered as a modification of the lowest order diagram shown in Fig. 7.1 by the substitution of one of the three loop diagrams depicted in Fig. 7.4. Diagram 7.4 (a) is the electron self-energy diagram, Diagram 7.4 (b) corresponds to the photon self-energy which is commonly referred to as the vacuum polarization diagram and Diagram 7.4 (c) corresponds to the modification of the basic vertex diagram representing Fermion-photon interaction. Each of the diagrams in Fig. 7.3


Figure 7.4: Feynman diagrams representing (a) Self energy of the electron (b) Self energy of the photon or Vacuum polarization and (c) Vertex correction.
uses one of the basic modifications depicted in Fig. 7.4.
Diagrams 7.3(a) and 7.3(b) include the basic self-energy diagram 7.4(a); Diagram 7.4(c) includes the basic photon self-energy diagram 7.4(b) and Diagram 7.4(d) includes the basic modification in the electron-photon vertex. It is possible to extract the effect of these loop diagrams depicted in Fig. 7.4 on higher order terms in perturbation theory. The loop diagrams lead to divergent integrals and it is a major problem to regularize these integrals and extract meaningful results from them. We have already seen how one can obtain meaningful results from the self-energy diagram.

$$
\begin{align*}
\mathcal{M}_{3}^{(a)} & =-4 \pi e^{3} \int\left(\bar{u}_{f} \not d \frac{1}{\not p_{i}-m} \gamma_{\mu} \frac{1}{\not p_{i}-\not p-m} \gamma_{\mu} u_{i}\right) \frac{1^{2}}{\mathbf{k}} \frac{d^{4} k}{(2 \pi)^{4}} .  \tag{7.21}\\
\mathcal{M}_{3}^{(b)} & =-4 \pi e^{3} \int\left(\bar{u}_{f} \gamma_{\mu} \frac{1}{\not p_{f}-\not x-m} \gamma_{\mu} \frac{1}{\not p_{f}-m} d u_{i}\right) \frac{1^{2}}{\mathbf{k}} \frac{d^{4} k}{(2 \pi)^{4}} .  \tag{7.22}\\
\mathcal{M}_{3}^{(c)} & =-4 \pi e^{3}\left(\bar{u}_{f} \gamma_{\mu} u_{1}\right) \frac{1}{\mathbf{q}^{2}} \int \sum_{\text {spins }}\left(\bar{u} \frac{1}{\not p-m} \gamma_{\mu} \frac{1}{p b+\not q-m} d u\right) \frac{d^{4} p}{(2 \pi)^{4}} ;  \tag{7.23}\\
\mathcal{M}_{3}^{(d)} & =-4 \pi e^{3} \int\left(\bar{u}_{f} \gamma_{\mu} \frac{1}{\not p_{f}-\not x-m} d \frac{1}{\not p_{i}-\nmid-m} \gamma_{\mu} u_{i}\right) \frac{1^{2}}{\mathbf{k}} \frac{d^{4} k}{(2 \pi)^{4}} . \tag{7.24}
\end{align*}
$$

Comparing the higher order matrix elements (7.21) - (7.24) with the lowest order matrix element (7.20), we can rewrite Eqs. (7.21) - (7.24) as given
below:

$$
\begin{align*}
\mathcal{M}_{3}^{(a)} & =-i e \bar{u}_{f} \not d \Sigma\left(\boldsymbol{p}_{i}\right) u_{i}  \tag{7.25}\\
\mathcal{M}_{3}^{(b)} & =-i e \bar{u}_{f} \Sigma\left(\boldsymbol{p}_{f}\right) d u_{i}  \tag{7.26}\\
\mathcal{M}_{3}^{(c)} & =-i e \bar{u}_{f} \gamma_{\mu} u_{i} \Pi  \tag{7.27}\\
\mathcal{M}_{3}^{(d)} & =-i e \bar{u}_{f} \gamma_{\mu} \Lambda u_{i} \tag{7.28}
\end{align*}
$$

where

$$
\begin{align*}
& \Sigma\left(\boldsymbol{p}_{i}\right)=-i 4 \pi e^{2} \int\left(\frac{1}{\not p_{i}-m} \gamma_{\mu} \frac{1}{\not p_{i}-\not \nsim-m} \gamma_{\mu} u_{i}\right) \frac{1^{2}}{\mathbf{k}} \frac{d^{4} k}{(2 \pi)^{4}} ;  \tag{7.29}\\
& \Sigma\left(\boldsymbol{p}_{f}\right)=-i 4 \pi e^{2} \int\left(\frac{1}{\not p_{f}-\not x-m} \gamma_{\mu} \frac{1}{\not p_{f}-m} \gamma_{\mu} u_{i}\right) \frac{1^{2}}{\mathbf{k}} \frac{d^{4} k}{(2 \pi)^{4}} ;  \tag{7.30}\\
& \Pi(\boldsymbol{q})=-i 4 \pi e^{2} \int \sum_{\text {spins }} \frac{1}{\mathbf{q}^{2}}\left(\bar{u} \frac{1}{\not p-m} \gamma_{\mu} \frac{1}{\not p+\not q-m} d u\right) \frac{d^{4} p}{(2 \pi)^{4}} ;  \tag{7.31}\\
& \Lambda\left(\boldsymbol{p}_{f}, \boldsymbol{p}_{i}\right)=-i 4 \pi e^{2} \int\left(\gamma_{\mu} \frac{1}{\not p_{f}-\not p-m} \not d^{\not p} \frac{1}{\not p_{i}-\not p-m} \gamma_{\mu}\right) \frac{1^{2}}{\mathbf{k}} \frac{d^{4} k}{(2 \pi)^{4}} . \tag{7.32}
\end{align*}
$$

Thus, for calculating the second order corrections to the elastic scattering of electron by a static external potential, we need to evaluate only the three loop integrals $\Sigma, \Pi, \Lambda$ that occurs in expressions (7.29) - (7.31) and depicted in Fig. 7.4. All the three integrals are found to be divergent for large values of the momentum variables of integration. The concepts of mass and charge renormalization enable one to extract unambiguously finite radiative corrections of order $\alpha$, expressed in terms of observed mass and charge. The importance of this analysis is to point out that the radiative corrections of lowest order to any process involve the same three divergent integrals and once we have coped up with these three integrals, the calculations of the second order radiative corrections to any process presents no difficulties, in principle.

## Review Questions

7.1 Draw a Feynman diagram for calculating the self energy of electron and show that it involves divergent integrals. Explain the Feynman method of coping with this integral and show that the electron self-energy is equivalent to a shift in mass. Obtain an expression for the mass shift.
7.2 Explain the concept of mass renormalizaion and explain how one can obtain the Lamb shift for a bound electron in hydrogen atom.
7.3 Consider the elastic scattering of an electron by a static external potential and draw Feynman diagrams for obtaining second-order radiative corrections for this process.

## Problems

7.1 Show that for two operators $A$ and $B$, whether commuting or non-commuting, the following expansions are true:

$$
\begin{aligned}
& \frac{1}{A-B}=\frac{1}{A}+\frac{1}{A} B \frac{1}{A}+\frac{1}{A} B \frac{1}{A} B \frac{1}{A}+\cdots \\
& \frac{1}{A+B}=\frac{1}{A}-\frac{1}{A} B \frac{1}{A}+\frac{1}{A} B \frac{1}{A} B \frac{1}{A}-\cdots
\end{aligned}
$$

7.2 Prove the following Feynman's formulae:

$$
\begin{aligned}
& \text { (1) } \frac{1}{a b}=\int_{0}^{1} \frac{d z}{[a z+b(1-z)]^{2}} \\
& \text { (2) } \frac{1}{a^{2} b}=\int_{0}^{1} \frac{2 z d z}{[a z+b(1-z)]^{3}} \\
& \text { (3) } \frac{1}{a b^{2}}=\int_{0}^{1} \frac{2(1-z) d z}{[a z+b(1-z)]^{3}}
\end{aligned}
$$

7.3 Establish the following relation

$$
\int_{-\infty}^{\infty} \frac{\left(1, k_{\sigma}\right) d^{4} k}{(2 \pi)^{4}\left(\mathbf{k}^{2}+i \epsilon-L\right)^{3}}=\frac{(1,0)}{32 \pi^{2} i L}
$$

which is given in a compact form. The notation $\left(1: k_{\sigma}\right)$ in the numerator means either 1 or $k_{\sigma}$ and correspondingly the notation (1:0) on the right hand side of the equation means 1 or 0 .
7.4 Given the following relation

$$
\int_{-\infty}^{\infty} \frac{\left(1, k_{\sigma}\right) d^{4} k}{(2 \pi)^{4}\left(\mathbf{k}^{2}+i \epsilon-L\right)^{3}}=\frac{(1,0)}{32 \pi^{2} i L}
$$

change the variable $\mathbf{k} \rightarrow \mathbf{k}-\mathbf{p}$ and obtain the relation given below.

$$
\int_{-\infty}^{\infty} \frac{\left(1 ; k_{\sigma}\right) d^{4} k}{(2 \pi)^{4}\left(\mathbf{k}^{2}-2 \mathbf{p} \cdot \mathbf{k}-\Delta\right)^{3}}=\frac{\left(1 ; p_{\sigma}\right)}{32 \pi^{2} i\left(\mathbf{p}^{2}+\Delta\right)}
$$

7.5 Evaluate the following definite integrals and show that they yield the results given below:
(a) $\int_{0}^{1} \ln \left[x^{-2}(1-x)\right] d x=1$;
(b) $\int_{0}^{1} x \ln \left[x^{-2}(1-x)\right] d x=-1 / 4$.

## Solutions to Problems

7.1 Post-multiply the first of equations by $A-B$.

$$
\begin{align*}
\frac{1}{A-B}(A-B)= & \frac{1}{A}(A-B)+\frac{1}{A} B \frac{1}{A}(A-B)+\frac{1}{A} B \frac{1}{A} B \frac{1}{A}(A-B) \\
& +\cdots \\
1= & C_{1}+C_{2}+C_{3}+\cdots \tag{7.33}
\end{align*}
$$

where

$$
\begin{aligned}
C_{1} & =1-\frac{1}{A} B \\
C_{2} & =\frac{1}{A} B \frac{1}{A}(A-B) \\
& =\frac{1}{A} B\left(1-\frac{1}{A} B\right) \\
& =\frac{1}{A} B-\frac{1}{A} B \frac{1}{A} B \\
C_{3} & =\frac{1}{A} B \frac{1}{A} B \frac{1}{A}(A-B) \\
& =\frac{1}{A} B \frac{1}{A} B-\frac{1}{A} B \frac{1}{A} B \frac{1}{A} B
\end{aligned}
$$

On the right-hand side of Eq. (7.33), only the first term (unity) remains and all the succeeding terms mutually cancel away.
Thus we have proved the first relation. Just change the sign of $B$ in the relation 1 and we get the second relation.
7.2 Let us start with a simple identity.

$$
\begin{equation*}
\frac{1}{a b}=\frac{1}{b-a}\left(\frac{1}{a}-\frac{1}{b}\right)=\frac{1}{b-a} \int_{a}^{b} \frac{d x}{x^{2}} \tag{7.34}
\end{equation*}
$$

Let $x=a z+b(1-z)$, then $d x=(a-b) d z$. Substituting them on the right hand side of Eq. (7.34), we get

$$
\begin{equation*}
\frac{1}{a b}=\frac{1}{b-a} \int \frac{(a-b) d z}{[a z+b(1-z)]^{2}} \tag{7.35}
\end{equation*}
$$

The lower limit of integration $x=a$ corresponds to $z=1$ and the upper limit $x=b$ corresponds to $z=0$. Hence we get the relation (1)

$$
\frac{1}{a b}=\int_{0}^{1} \frac{d z}{[a z+b(1-z)]^{2}}
$$

Differentiating both sides of the above equation with respect to $a$, we get the relation (2).

$$
\frac{1}{a^{2} b}=\int_{0}^{1} \frac{2 z d z}{[a z+b(1-z)]^{3}}
$$

Differentiating both sides of Eq. (1) with respect to $b$, we get the relation (3)

$$
\frac{1}{a b^{2}}=\int_{0}^{1} \frac{2(1-z) d z}{[a z+b(1-z)]^{3}}
$$

7.3 First let us consider $k_{\sigma}$ in the numerator of the integral. Since $k_{\sigma}$ is a odd function, the integral vanishes and it corresponds to 0 on the right hand side of the equation.
To evaluate the integral with 1 in the numerator, we need to use the contour integration. Since $d^{4} k=d \omega d^{3} k$ and $\mathbf{k}^{2}=\omega^{2}-\boldsymbol{k}^{2}$, the factor in the denominator can be written as

$$
\mathbf{k}^{2}+i \epsilon-L=\omega^{2}+i \epsilon-\left(L+\boldsymbol{k}^{2}\right)
$$

Then for $\epsilon \ll L+\boldsymbol{k}^{2}$, there are poles at $\omega= \pm\left[\left(L+\boldsymbol{k}^{2}\right)^{1 / 2}-i \epsilon\right]$, and contour integration of $\omega$ yields

$$
\begin{equation*}
\int_{-\infty}^{\infty} \frac{d \omega}{\omega^{2}+i \epsilon-\left(L+\boldsymbol{k}^{2}\right)}=2 \pi i\left\{\frac{1}{-2\left(L+\boldsymbol{k}^{2}\right)^{1 / 2}}\right\} \tag{7.36}
\end{equation*}
$$

with the contour in the upper half-plane.
Let us differentiate both sides of Eq. (7.36) with respect to $L$. First let us differentiate the L.H.S. of Eq. (7.36).

$$
\begin{aligned}
& \frac{d}{d L}\left(\frac{1}{\omega^{2}+i \epsilon-\left(L+\boldsymbol{k}^{2}\right)}\right)=\frac{d}{d L}\left(\frac{1}{a-L}\right)=\frac{1}{(a-L)^{2}} \\
& \frac{d^{2}}{d L^{2}}\left(\frac{1}{a-L}\right)=\frac{d}{d L}\left(\frac{1}{(a-L)^{2}}\right)=\frac{2}{(a-L)^{3}}
\end{aligned}
$$

Now let us differentiate twice the R.H.S. of Eq. (7.36).

$$
\begin{aligned}
-\pi i \frac{d}{d L}\left(L+\boldsymbol{k}^{2}\right)^{-1 / 2} & =\frac{\pi i}{2}\left(L+\boldsymbol{k}^{2}\right)^{-3 / 2} \\
-\pi i \frac{d^{2}}{d L^{2}}\left(L+\boldsymbol{k}^{2}\right)^{-1 / 2} & =\frac{\pi i}{2} \frac{d}{d L}\left(L+\boldsymbol{k}^{2}\right)^{-3 / 2}=-\frac{3 \pi i}{4}\left(L+\boldsymbol{k}^{2}\right)^{-5 / 2}
\end{aligned}
$$

Equating the differentiated quantities on both sides, we get

$$
\begin{equation*}
\int \frac{d \omega}{\left\{\omega^{2}+i \epsilon-\left(L+\boldsymbol{k}^{2}\right)\right\}^{3}}=\frac{3 \pi}{8 i}\left(L+\boldsymbol{k}^{2}\right)^{-5 / 2} \tag{7.37}
\end{equation*}
$$

Of the four dimensional integration $d^{4} k$, we have performed the integration over the energy $d \omega$ and obtained the result (7.37). Now we have to perform the integration over three-momenta $d^{3} k$.

$$
\begin{aligned}
\int\left(L+\boldsymbol{k}^{2}\right)^{-5 / 2} d^{3} k & =4 \pi \int_{0}^{\infty}\left(L+\boldsymbol{k}^{2}\right)^{-5 / 2} k^{2} d k \\
& =\left.4 \pi \frac{k^{3}}{3 L\left(L+k^{2}\right)^{3 / 2}}\right|_{0} ^{\infty}=\frac{4 \pi}{3 L}
\end{aligned}
$$

Collecting all the results, we get finally

$$
\begin{align*}
\int \frac{d^{4} k}{(2 \pi)^{4}\left(\mathbf{k}^{2}+i \epsilon-L\right)^{3}} & =\int \frac{d \omega d^{3} k}{(2 \pi)^{4}\left\{\omega^{2}+i \epsilon-\left(L+\boldsymbol{k}^{2}\right)\right\}^{3}} \\
& =\frac{1}{(2 \pi)^{4}} \int \frac{3 \pi}{8 i}\left(L+\boldsymbol{k}^{2}\right)^{-5 / 2} d^{3} k \\
& =\frac{1}{(2 \pi)^{4}} \frac{3 \pi}{8 i} \frac{4 \pi}{3 L}=\frac{1}{32 \pi^{2} i L} \tag{7.38}
\end{align*}
$$

7.4 Let us start from the relation

$$
\int_{-\infty}^{\infty} \frac{\left(1, k_{\sigma}\right) d^{4} k}{(2 \pi)^{4}\left(\mathbf{k}^{2}+i \epsilon-L\right)^{3}}=\frac{(1,0)}{32 \pi^{2} i L}
$$

Change the variable $\mathbf{k} \rightarrow \mathbf{k}^{\prime}=\mathbf{k}-\mathbf{p}$. Then $\mathbf{k}^{\prime 2}=\mathbf{k}^{2}+\mathbf{p}^{2}-2 \mathbf{p} \cdot \mathbf{k}$.

$$
\begin{align*}
\int \frac{d^{4} k^{\prime}}{(2 \pi)^{4}} \frac{1}{\left(\mathbf{k}^{\prime 2}+i \epsilon-L\right)^{3}} & =\int \frac{d^{4} k^{\prime}}{(2 \pi)^{4}} \frac{1}{\left\{\mathbf{k}^{2}-2 \mathbf{p} \cdot \mathbf{k}+i \epsilon-\left(L-\mathbf{p}^{2}\right)\right\}^{3}} \\
= & \int \frac{d^{4} k^{\prime}}{(2 \pi)^{4}} \frac{1}{\left\{\mathbf{k}^{2}-2 \mathbf{p} \cdot \mathbf{k}+i \epsilon-\Delta\right\}^{3}}  \tag{7.39}\\
& \text { where } \Delta=L-\mathbf{p}^{2} \\
= & \int \frac{d^{4} k^{\prime}}{(2 \pi)^{4}} \frac{1}{\left\{\mathbf{k}^{\prime 2}+i \epsilon-\left(\Delta+\mathbf{p}^{2}\right)\right\}^{3}} \\
= & \frac{1}{32 \pi^{2} i\left(\Delta+\mathbf{p}^{2}\right)} \tag{7.40}
\end{align*}
$$

Now, let us consider the other integrand.

$$
\int \frac{d^{4} k^{\prime}}{(2 \pi)^{4}} \frac{k_{\sigma}}{\left\{\mathbf{k}^{\prime 2}+i \epsilon-\left(\Delta+\mathbf{p}^{2}\right)\right\}^{3}}=\int \frac{d^{4} k^{\prime}}{(2 \pi)^{4}} \frac{k_{\sigma}^{\prime}+p_{\sigma}}{\left\{\mathbf{k}^{\prime 2}+i \epsilon-\left(\Delta+\mathbf{p}^{2}\right)\right\}^{3}}
$$

The term $k_{\sigma}^{\prime}$ will yield zero since it is an odd function in $\mathbf{k}$. So, the term $p_{\sigma}$ alone will contribute to the integral. Hence

$$
\begin{equation*}
\int \frac{d^{4} k^{\prime}}{(2 \pi)^{4}} \frac{k_{\sigma}^{\prime}+p_{\sigma}}{\left\{\mathbf{k}^{\prime 2}+i \epsilon-\left(\Delta+\mathbf{p}^{2}\right)\right\}^{3}}=\frac{p_{\sigma}}{32 \pi^{2} i\left(\Delta+\mathbf{p}^{2}\right)} \tag{7.41}
\end{equation*}
$$

Combining Eqs. (7.40) and (7.41) and changing the dummy variable $d^{4} k^{\prime}$ to $d^{4} k$, we get

$$
\begin{equation*}
\int \frac{d^{4} k}{(2 \pi)^{4}} \frac{1 ; k_{\sigma}+p_{\sigma}}{\left\{\mathbf{k}^{2}+i \epsilon-\left(\Delta+\mathbf{p}^{2}\right)\right\}^{3}}=\frac{1 ; p_{\sigma}}{32 \pi^{2} i\left(\Delta+\mathbf{p}^{2}\right)} \tag{7.42}
\end{equation*}
$$

7.5 To evaluate the given definite integrals, we need to use the following standard integrals.

$$
\int \ln x d x=x \ln x-x ; \quad \int x \ln x d x=\frac{x^{2}}{2} \ln x-\frac{x^{2}}{4}
$$

One can easily check the above relations by differentiating the R.H.S. and checking whether it yields the integrand on the L.H.S.
(a) Let us first expand $\ln \left(\frac{1-x}{x^{2}}\right)$.

$$
\ln \left(\frac{1-x}{x^{2}}\right)=\ln (1-x)-\ln x^{2}=\ln (1-x)-2 \ln x
$$

Then

$$
\begin{aligned}
\int_{0}^{1} \ln \left(\frac{1-x}{x^{2}}\right) d x= & \int_{0}^{1} \ln (1-x) d x-2 \int_{0}^{1} \ln x d x \\
\int_{0}^{1} \ln (1-x) d x= & -\int_{1}^{0} \ln y d y \\
& \text { substituting } y=1-x \text { and } d y=-d x \\
= & -|y \ln y-y|_{1}^{0}=-1 \\
\int_{0}^{1} \ln x d x= & |x \ln x-x|_{0}^{1}=-1
\end{aligned}
$$

Collecting the above results, we get

$$
\int_{0}^{1} \ln \left(\frac{1-x}{x^{2}}\right) d x=1
$$

(b) A similar procedure is adopted.

$$
\begin{aligned}
\int_{0}^{1} x \ln \left(\frac{1-x}{x^{2}}\right) d x & =\int\left\{x \ln (1-x)-x \ln x^{2}\right\} d x \\
& =\int x \ln (1-x) d x-2 \int x \ln x d x \\
\int_{0}^{1} x \ln (1-x) d x & =-\int(1-y) \ln y d y \\
& \text { substituting } y=1-x \text { and } d y=-d x \\
& =-\int_{1}^{0} \ln y d y+\int_{1}^{0} y \ln y d y \\
& =-|y \ln y-y|_{1}^{0}+\left|\left(x^{2} / 2\right) \ln x-\left(x^{2} / 4\right)\right|_{0}^{1} \\
& =-1+\frac{1}{4}=-\frac{3}{4} \\
\int_{0}^{1} x \ln x d x & =\left|\frac{x^{2}}{2} \ln x-\frac{x^{2}}{4}\right|_{0}^{1}=-\frac{1}{4} .
\end{aligned}
$$

Collecting the above results, we get

$$
\int_{0}^{1} x \ln \left(\frac{1-x}{x^{2}}\right) d x=-\frac{1}{4}
$$

## Chapter 8

## Elements of Quantum Field Theory

The formulation of Quantum Mechanics is based on the Hamiltonian theory of classical mechanics by treating the canonically conjugate dynamical variables, position and momentum as operators, obeying certain commutation relations. This is known as the first quantization. The alternative Lagrangian formulation of classical mechanics and the Hamiltonian action principle can be considered more fundamental than the Newtonian dynamics and they lead to the formulation of both classical and quantum field theory. The Lagrangian method can be extended to relativistic fields since the action function is relativistically invariant.

A field is specified by a number of functions, say $n$, of space-time.

$$
\psi_{\rho}\left(x_{1}, x_{2}, x_{3}, t\right), \quad \rho=1,2,3, \ldots, n
$$

The classic example is that of electromagnetic field, which is described by the six components of the electric and magnetic field strengths at each point in space-time. A great merit of field theory is that the special theory of relativity can be incorporated in it in a simple way. Let us replace the time variable $t$ by

$$
x_{0}=c t
$$

so that all the coordinates $x_{0}, x_{1}, x_{2}, x_{3}$ now have the dimension of length. The dynamical character of the field will be described by the field equations. They will be assumed to be partial differential equations of order, not greater than two. It will be further assumed that they can be derived
by a variational principle, in the same way as the Lagrangian equations of motion for particles are obtained from the action principle.

Let us start with a brief review of classical mechanics ${ }^{1}$ since most of the concepts of field theory have been drawn from the LagrangianHamiltonian formalism.

### 8.1 A brief review of classical mechanics

### 8.1.1 The Lagrangian-Hamiltonian formalism

Let us recapitulate the derivation of the Lagrange equations of motion from the action principle. Let $q_{k}(t), k=1,2, \ldots, n$ be the generalized coordinates and $\dot{q}_{k}(t)=\frac{d q_{k}(t)}{d t}$, the generalized velocities of a mechanical system and $L$, its Lagrangian.

$$
L=L\left(q_{k}(t), \dot{q}_{k}(t), t\right)
$$

The action integral is

$$
\begin{equation*}
\mathscr{A}=\int_{t_{1}}^{t_{2}} L\left(q_{k}, \dot{q}_{k}, t\right) d t \tag{8.1}
\end{equation*}
$$

Now consider the variation of $q_{k}(t)$,

$$
q_{k}(t) \longrightarrow q_{k}(t)+\delta q_{k}(t),
$$

where $\delta q_{k}(t)$ is an arbitrary infinitesimal change, which vanishes at the end-points $t_{1}$ and $t_{2}$.

$$
\delta q_{k}\left(t_{1}\right)=\delta q_{k}\left(t_{2}\right)=0
$$

This variation results in the change of $\dot{q}_{k}(t)$ to

$$
\dot{q}_{k}(t)+\frac{d}{d t} \delta q_{k}(t) .
$$

The variational principle ${ }^{2}$ asserts that the action integral has a stationary value and the variation $\delta \mathscr{A}$ of $\mathscr{A}$ due to small variation in path with fixed end-points is identically zero.

$$
\begin{equation*}
\delta \mathscr{A}=\int_{t_{1}}^{t_{2}} \sum_{k}\left\{\frac{\partial L}{\partial q_{k}} \delta q_{k}+\frac{\partial L}{\partial \dot{q}_{k}} \frac{d}{d t} \delta q_{k}\right\} d t=0 . \tag{8.2}
\end{equation*}
$$

[^56]The second term in the integral is evaluated using the principle of integration by parts.

$$
\begin{align*}
\int \frac{\partial L}{\partial \dot{q}_{k}} \frac{d}{d t} \delta q_{k} d t & =\left.\frac{\partial L}{\partial \dot{q}_{k}} \delta q_{k}\right|_{t_{1}} ^{t_{2}}-\int \frac{d}{d t} \frac{\partial L}{\partial \dot{q}_{k}} \delta q_{k} d t \\
& =-\int \frac{d}{d t} \frac{\partial L}{\partial \dot{q}_{k}} \delta q_{k} d t \tag{8.3}
\end{align*}
$$

since $\left.\frac{\partial L}{\partial \dot{q}_{k}} \delta q_{k}\right|_{t_{1}} ^{t_{2}}$ vanishes because $\delta q_{k}\left(t_{1}\right)=\delta q_{k}\left(t_{2}\right)=0$, due to the fixed end-points.

Substituting (8.3) into Eq. (8.2), we get

$$
\begin{equation*}
\delta \mathscr{A}=\int_{t_{1}}^{t_{2}} \sum_{k}\left\{\frac{\partial L}{\partial q_{k}}-\frac{d}{d t} \frac{\partial L}{\partial \dot{q}_{k}}\right\} \delta q_{k} d t=0 . \tag{8.4}
\end{equation*}
$$

Since the generalized coordinates $q_{k}$ are independent and the variations $\delta q_{k}$ are arbitrary, the condition that $\delta \mathscr{A}=0$ requires that the coefficients of $\delta q_{k}$ in the integrand of (8.4) should separately vanish.

$$
\begin{equation*}
\frac{\partial L}{\partial q_{k}}-\frac{d}{d t} \frac{\partial L}{\partial \dot{q}_{k}}=0, \quad k=1,2, \ldots, n \tag{8.5}
\end{equation*}
$$

Equation (8.5) is known as the Euler-Lagrange equation of motion.
The momentum $p_{k}$ conjugate to the coordinate $q_{k}$ is given by

$$
\begin{equation*}
p_{k}=\frac{\partial L}{\partial \dot{q}_{k}} . \tag{8.6}
\end{equation*}
$$

Differentiating $p_{k}$ with respect to $t$ and using the relation (8.5), we obtain

$$
\begin{equation*}
\dot{p}_{k}=\frac{\partial L}{\partial q_{k}} . \tag{8.7}
\end{equation*}
$$

If $T$ is the kinetic energy and $V$ is the potential energy, then the Lagrangian of the system is given by

$$
\begin{equation*}
L=T-V, \tag{8.8}
\end{equation*}
$$

and the Hamiltonian function is given by

$$
\begin{equation*}
H\left(q_{k}, p_{k}, t\right)=\sum_{k} p_{k} \dot{q}_{k}-L . \tag{8.9}
\end{equation*}
$$

It follows that

$$
\frac{\partial H}{\partial p_{k}}=\dot{q}_{k} \quad \text { and } \quad \frac{\partial H}{\partial q_{k}}=-\frac{\partial L}{\partial q_{k}} .
$$

Using the Euler-Lagrange Eq. (8.5) and the relation (8.6), we deduce Hamilton's canonical equations of motion.

$$
\begin{equation*}
\dot{q}_{k}=\frac{\partial H}{\partial p_{k}} ; \quad \dot{p}_{k}=-\frac{\partial H}{\partial q_{k}} . \tag{8.10}
\end{equation*}
$$

From (8.10), one can deduce the equations of motion for a general dynamical variable $F \equiv F\left(q_{k}, p_{k}, t\right)$.

$$
\begin{align*}
\frac{d F}{d t} & =\frac{\partial F}{\partial t}+\sum_{k}\left(\frac{\partial F}{\partial q_{k}} \frac{\partial q_{k}}{\partial t}+\frac{\partial F}{\partial p_{k}} \frac{\partial p_{k}}{\partial t}\right) \\
& =\frac{\partial F}{\partial t}+[F, H]_{\mathrm{PB}}, \tag{8.11}
\end{align*}
$$

where

$$
\begin{equation*}
[F, H]_{\mathrm{PB}}=\sum_{k}\left(\frac{\partial F}{\partial q_{k}} \frac{\partial H}{\partial p_{k}}-\frac{\partial H}{\partial q_{k}} \frac{\partial F}{\partial p_{k}}\right) \tag{8.12}
\end{equation*}
$$

is the Poisson Bracket.

### 8.1.2 The classical fields

From a discrete mechanical system, it is possible to go to a continuous system when the number of degrees of freedom becomes infinite. Consider a classical field $\psi(\boldsymbol{x}, t)$ which is continuous and described by an infinite set of parameters $\boldsymbol{x}$ and $t$. The Lagrangian formalism can be extended to the fields and the variational principle can be incorporated in an analogous manner by introducing certain modifications.

1. Since the field $\psi(\boldsymbol{x}, t)$ depends on the continuous variable $\boldsymbol{x}$, we have to invoke the concept of Lagrangian density.
2. The theory can be made relativistically invariant by treating both $\boldsymbol{x}$ and $t$ coordinates on equal footing. This can be done by using various metrics but we find it more convenient to use the metric
$1,-1,-1,-1$ instead of the Minkowski space ${ }^{3}$, with the time coordinate $x_{0}=c t$. By this procedure, we make all the four coordinates of the same dimension of length.

Let us start with the Lagrangian density $\mathscr{L}$.

$$
\mathscr{L}\left(\psi_{\rho}, \frac{\partial \psi_{\rho}}{\partial x_{\mu}}, \mathbf{x}\right) \equiv \mathscr{L}\left(\psi_{\rho}, \psi_{\rho, \mu}, \mathbf{x}\right)
$$

where, for brevity, a notation $\psi_{\rho, \mu}$ is used to denote $\frac{\partial \psi_{\rho}}{\partial x_{\mu}}$. The Lagrangian density is a function of field functions $\psi_{\rho}$, its first order derivatives and the space-time coordinates ( $\mathrm{x}: x_{0}=c t, x_{1}, x_{2}, x_{3}$ ). It is important that $\mathscr{L}$ should not depend on second and higher order derivatives of $\psi_{\rho}$.

The Lagrangian and the action integral are given by

$$
\begin{align*}
L & =\int \mathscr{L}\left(\psi_{\rho}, \psi_{\rho, \mu}, \mathbf{x}\right) d^{3} x  \tag{8.13}\\
\mathscr{A} & =\int_{V} \mathscr{L}\left(\psi_{\rho}, \psi_{\rho, \mu}, \mathbf{x}\right) d^{4} x \tag{8.14}
\end{align*}
$$

where

$$
d^{4} x=d x_{0} d x_{1} d x_{2} d x_{3}=c d t d x d y d z
$$

and $V$ is a certain domain in space-time, bounded by a hypersurface $S$. According to the principle of least action, the variation $\delta \mathscr{A}$ in $\mathscr{A}$ for arbitrary domains $V$ is zero for the variations of $\psi_{\rho}(\mathbf{x})$.

$$
\psi_{\rho}(\mathbf{x}) \longrightarrow \psi_{\rho}(\mathbf{x})+\delta \psi_{\rho}(\mathbf{x})
$$

where $\delta \psi_{\rho}(\mathbf{x})$ is an arbitrary infinitesimal variation of the first order, which is continuous and differentiable and vanishes on the boundary $S$ of $V$, i.e..

$$
\begin{equation*}
\delta \psi_{\rho}(\mathbf{x})=0, \quad \text { on the surface } S \text { of } V \tag{8.15}
\end{equation*}
$$

[^57]This causes a change in the derivative.

$$
\frac{\partial \psi_{\rho}}{\partial x_{\mu}} \longrightarrow \frac{\partial \psi_{\rho}}{\partial x_{\mu}}+\frac{\partial}{\partial x_{\mu}} \delta \psi_{\rho} .
$$

Thus

$$
\begin{equation*}
\delta \mathscr{A}=\int_{V} \sum_{\rho}\left\{\frac{\partial \mathscr{L}}{\partial \psi_{\rho}} \delta \psi_{\rho}+\sum_{\mu} \frac{\partial \mathscr{L}}{\partial \frac{\partial \psi_{\rho}}{\partial x_{\mu}}} \frac{\partial}{\partial x_{\mu}} \delta \psi_{\rho}\right\} d^{4} x=0 . \tag{8.16}
\end{equation*}
$$

Integrating the second term in the curly brackets by parts, we get

$$
\begin{equation*}
\int_{V} \frac{\partial \mathscr{L}}{\partial \frac{\partial \psi_{\rho}}{\partial x_{\mu}}} \frac{\partial}{\partial x_{\mu}} \delta \psi_{\rho} d^{4} x=\left.\frac{\partial \mathscr{L}}{\partial \frac{\partial \psi_{\rho}}{\partial x_{\mu}}} \delta \psi_{\rho}\right|_{S}-\int_{V} \frac{\partial}{\partial x_{\mu}} \frac{\partial \mathscr{L}}{\partial \frac{\partial \psi_{\rho}}{\partial x_{\mu}}} \delta \psi_{\rho} d^{4} x . \tag{8.17}
\end{equation*}
$$

The first term on the right hand side of Eq. (8.17) vanishes because of the boundary condition (8.15) and hence the variational principle leads to the result

$$
\begin{equation*}
\delta \mathscr{A}=\int_{V} \sum_{\rho}\left\{\frac{\partial \mathscr{L}}{\partial \psi_{\rho}}-\sum_{\mu} \frac{\partial}{\partial x_{\mu}} \frac{\partial \mathscr{L}}{\partial \frac{\partial \psi_{\rho}}{\partial x_{\mu}}}\right\} \delta \psi_{\rho} d^{4} x=0 . \tag{8.18}
\end{equation*}
$$

Since this equation holds for arbitrary $\delta \psi_{\rho}$ and also valid for arbitrary domain $V$, the resulting relations

$$
\begin{equation*}
\frac{\partial \mathscr{L}}{\partial \psi_{\rho}}-\sum_{\mu} \frac{\partial}{\partial x_{\mu}} \frac{\partial \mathscr{L}}{\partial \frac{\partial \psi_{\rho}}{\partial x_{\mu}}}=0, \quad \rho=1,2, \ldots, n \tag{8.19}
\end{equation*}
$$

which are analogous to the Euler-Lagrange equations (8.5) are obtained. These equations, in turn, lead to the field equations. Since $\mathscr{L}$ does not involve derivatives of $\psi_{\rho}$ of order higher than the first, the resulting field equations will be utmost of second order.

One could, in principle, consider field equations that are not derivable from a variational principle, just as in mechanics we can have systems with frictional and dissipative forces for which equations of motion cannot be so obtained. Non-conservative forces, like friction, appear at a macroscopic level since we wish to neglect microscopic complications. So, at the microscopic level, the fundamental laws can be put in the form of the principle of least action.

Using the short-hand notation, the Euler-Lagrange Eq. (8.19) can be rewritten as

$$
\begin{equation*}
\frac{\partial \mathscr{L}}{\partial \psi_{\rho}}-\partial_{\mu} \frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \psi_{\rho}\right)}=0, \quad \rho=1,2, \ldots, n \tag{8.20}
\end{equation*}
$$

with the convention that a summation is to be made on the repeated indices.

The conjugate momentum density $\Pi(\mathbf{x})$ and the Hamiltonian density $\mathscr{H}(\mathbf{x})$ are obtained in a way, analogous to Eqs. (8.6) and (8.9).

$$
\begin{align*}
\Pi(\mathbf{x}) & =\frac{\partial \mathscr{L}}{\partial \dot{\psi}} .  \tag{8.21}\\
\mathscr{H}(\mathbf{x}) & =\Pi(\mathbf{x}) \dot{\psi}(\mathbf{x})-\mathscr{L} . \tag{8.22}
\end{align*}
$$

### 8.2 Quantization of the field

The transition from the classical field theory to the quantum field theory is made by postulating the field variables $\psi$ and $\Pi$ as field operators and prescribing certain algebraic relations between them. The prescription of these algebraic relations (commutators or anticommutators) are known as the second quantization.

It may be recalled that some of the field equations, for example, the Schrödinger equation, the Klein-Gordon equation and the Dirac equation were obtained earlier by treating the dynamical variables $\boldsymbol{x}$ and $\boldsymbol{p}$ as operators obeying certain commutation relations. That is known as the first quantization. Since in quantum field theory, the field functions are treated as operators and certain algebraic relations are postulated between them for the second time, this is known as the second quantization.

The quantization of the field results in the description of the field in terms of "particles" or more precisely, "the field quanta" which are the carriers of energy, momentum and charge of the field. Let us illustrate the method of quantization by considering the Schrödinger field. It is a non-relativistic field and it allows the quantization procedure by choosing either the commutation relations or anticommutation relations between the field operators. Later, we shall see that in the case of relativistic fields such as the Klein-Gordon field or the Dirac field, only one of them is admissible.

### 8.2.1 The Schrödinger field

Let us consider the Schrödinger field which satisfies the field equation

$$
\begin{equation*}
i \hbar \frac{d \psi}{d t}+\frac{\hbar^{2}}{2 m} \nabla^{2} \psi-V \psi=0 . \tag{8.23}
\end{equation*}
$$

Quantum mechanics treats Eq. (8.23) as the equation of motion of a particle of mass $m$ moving in an external potential $V$. Here, we treat it as a classical field equation which, when quantized, will yield an assembly of a large number indistinguishable particles. Since this is the second time that quantization is done, this is known as the second quantization. It may be recalled that Eq. (8.23) has been obtained earlier by treating the dynamical variables $\boldsymbol{x}$ and $\boldsymbol{p}$ as operators obeying the commutation relation $\left[x, p_{x}\right]_{-}=i \hbar$; which is known as the first quantization.

The Lagrangian density for the Schrödinger field is given by

$$
\begin{equation*}
\mathscr{L}_{S}=i \hbar \psi^{*} \dot{\psi}-\frac{\hbar^{2}}{2 m} \nabla \psi^{*} \cdot \nabla \psi-V \psi^{*} \psi \tag{8.24}
\end{equation*}
$$

which, when substituted in the Euler-Lagrange Equation

$$
\begin{equation*}
\frac{\partial \mathscr{L}}{\partial \psi}-\nabla \cdot \frac{\partial \mathscr{L}}{\partial(\nabla \psi)}-\frac{\partial}{\partial t}\left(\frac{\partial \mathscr{L}}{\partial \dot{\psi}}\right)=0 \tag{8.25}
\end{equation*}
$$

yields Eq. (8.23). Equation (8.25) is the same as Eq. (8.20) but given in the expanded form, separating the space and time components.

The conjugate field $\Pi(\boldsymbol{x}, t)$ is given by

$$
\begin{equation*}
\Pi(\boldsymbol{x}, t)=\frac{\partial \mathscr{L}_{S}}{\partial \dot{\psi}}=i \hbar \psi^{*}(\boldsymbol{x}, t) \tag{8.26}
\end{equation*}
$$

The Hamiltonian density $\mathscr{H}$ and the Hamiltonian $H$ are given by

$$
\begin{align*}
\mathscr{H} & =\Pi \dot{\psi}-\mathscr{L}_{S}=\frac{\hbar^{2}}{2 m} \boldsymbol{\nabla} \psi^{*} \cdot \boldsymbol{\nabla} \psi-V \psi^{*} \psi:  \tag{8.27}\\
H & =\int_{V} \mathscr{H} d^{3} x=\int_{V}\left(\frac{\hbar^{2}}{2 m} \boldsymbol{\nabla} \psi^{*} \cdot \boldsymbol{\nabla} \psi-V \psi^{*} \psi\right) d^{3} x \tag{8.28}
\end{align*}
$$

where the suffix $V$ on the integral denotes the volume of integration, over which the field extends.

## Quantization

The quantization of the Schrödinger field can be done by expanding the field operator $\psi$ in terms of the complete set of eigenfunctions $u_{k}(\boldsymbol{x})$ of the field equation (8.23). This is known as the Fourier decomposition of the field into its normal modes.

$$
\begin{equation*}
\psi(\boldsymbol{x}, t)=\sum_{k} a_{k}(t) u_{k}(\boldsymbol{x}) \tag{8.29}
\end{equation*}
$$

where $a_{k}(t)$ is the Fourier coefficient which is interpreted as the annihilation operator of the field quantum with momentum $\boldsymbol{k}$. The functions $u_{k}(\boldsymbol{x})$ satisfy the following relations:

$$
\begin{align*}
\int u_{k}^{*}(\boldsymbol{x}) u_{l}(\boldsymbol{x}) d^{3} x & =\delta_{k l} ;  \tag{8.30}\\
\sum_{k} u_{k}(\boldsymbol{x}) u_{k}^{*}\left(\boldsymbol{x}^{\prime}\right) & =\delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right) \tag{8.31}
\end{align*}
$$

and they are the single particle eigenfunctions of the Hamiltonian

$$
H_{P}=-\frac{\hbar^{2}}{2 m} \nabla^{2}+V
$$

with eigenvalue $\epsilon_{k}$.

$$
H_{P} u_{k}(\boldsymbol{x})=\epsilon_{k} u_{k}(\boldsymbol{x}) .
$$

The Hermitian conjugate of $\psi$ is

$$
\psi^{\dagger}(\boldsymbol{x}, t)=\sum_{k} a_{k}^{\dagger}(t) u_{k}^{*}(\boldsymbol{x}) .
$$

The quantization is done by postulating certain algebraic relations between the field operators $\psi$ and $\psi^{\dagger}$ or alternatively between the operators $a_{k}(t)$ and $a_{l}^{\dagger}(t)$.

In case of the Schrödinger field, which is a non-relativistic field, it is possible to choose commutation relations or anticommutation relations ${ }^{4}$ between $a_{k}$ and $a_{k}^{\dagger}$. The first one corresponds to a system of Bosons and the second one corresponds to a system of Fermions.

[^58]
### 8.2.2 Quantization into Bosons

Let us postulate the following commutation rules for the field operator $\psi(\boldsymbol{x}, t)$ and its conjugate momentum $\Pi(\boldsymbol{x}, t)$ defined by Eq. (8.26):

$$
\begin{align*}
& {\left[\psi(\boldsymbol{x}, t), \Pi\left(\boldsymbol{x}^{\prime}, t\right)\right]_{-}=i \hbar \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)}  \tag{8.32}\\
& {\left[\psi(\boldsymbol{x}, t), \psi\left(\boldsymbol{x}^{\prime}, t\right)\right]_{-}=0=\left[\Pi(\boldsymbol{x}, t), \Pi\left(\boldsymbol{x}^{\prime}, t\right)\right]_{-}} \tag{8.33}
\end{align*}
$$

This is analogous to the Heisenberg commutation relations used in Quantum Mechanics for $\boldsymbol{x}$ and $\boldsymbol{p}$. This, in turn, leads to the commutation relations between the Fourier coefficients $a_{k}, a_{k}^{\dagger}$ on using the expansion (8.29).

$$
\begin{align*}
{\left[a_{k}, a_{l}^{\dagger}\right]_{-} } & =\delta_{k l} ;  \tag{8.34}\\
{\left[a_{k}, a_{l}\right]_{-} } & =\left[a_{k}^{\dagger}, a_{l}^{\dagger}\right]_{-}=0 \tag{8.35}
\end{align*}
$$

where all the operators refer to the same time. The operators $a_{k}$ and $a_{k}^{\dagger}$ are called the annihilation and creation operators of the field quanta. The entire description of the system - its state vector, the number and energy of the system and the effects of the operators on the system - can hereafter be studied in terms of these operators.

1. The number operator $N_{k}$ is an Hermitian operator

$$
\begin{equation*}
N_{k}=a_{k}^{\dagger} a_{k} \tag{8.36}
\end{equation*}
$$

with eigenvalues $n_{k}=0,1, \ldots, \infty$.
2. The state vector is better described in the occupation number representation, better known as the Fock space, giving the number of quanta in each state.

$$
\left|n_{1}, n_{2}, \ldots, n_{k}, \ldots\right\rangle=c\left(a_{1}^{\dagger}\right)^{n_{1}}\left(a_{2}^{\dagger}\right)^{n_{2}} \ldots\left(a_{k}^{\dagger}\right)^{n_{k}} \ldots|0\rangle,
$$

with

$$
c=\frac{1}{\left(n_{1}!n_{2}!\ldots n_{k}!\ldots\right)^{1 / 2}}
$$

The vacuum state $|0\rangle$ is defined by

$$
N_{k}|0\rangle=0, \quad \text { for all } k
$$

3. The effect of these operators on the state vector is given below:

$$
\begin{align*}
a_{k}\left|n_{1}, n_{2}, \ldots, n_{k}, \ldots\right\rangle & =\sqrt{n_{k}}\left|n_{1}, n_{2}, \ldots, n_{k}-1, \ldots\right\rangle ;  \tag{8.37}\\
a_{k}^{\dagger}\left|n_{1}, n_{2}, \ldots, n_{k}, \ldots\right\rangle & =\sqrt{n_{k}+1}\left|n_{1}, n_{2}, \ldots, n_{k}+1, \ldots\right\rangle ;  \tag{8.38}\\
N_{k}\left|n_{1}, n_{2}, \ldots, n_{k}, \ldots\right\rangle & =n_{k}\left|n_{1}, n_{2}, \ldots, n_{k}, \ldots\right\rangle . \tag{8.39}
\end{align*}
$$

4. The Hamiltonian of the system and its energy are given by

$$
\begin{align*}
H & =\sum_{k l} a_{k}^{\dagger} a_{l} \int_{V}\left(\frac{\hbar^{2}}{2 m} \nabla u_{k}^{*} \cdot \nabla u_{l}+V u_{k}^{*} u_{l}\right) d^{3} x  \tag{8.40}\\
& =\sum_{k} N_{k} \epsilon_{k} \tag{8.41}
\end{align*}
$$

The total energy of the field in the state $\left|n_{1}, n_{2}, \ldots, n_{k}, \ldots\right\rangle$ is given by

$$
E=\sum_{k} n_{k} \epsilon_{k} .
$$

Since the given state $u_{k}$ can be occupied by any number of particles, the field represents an assembly of Bosons.

### 8.2.3 Quantization into Fermions

For a system of Fermions, the exclusion principle restricts the occupation number $n_{k}$ of a particular state to 0 or 1 . It has been shown by Jordan and Wigner ${ }^{5}$ that this condition would be satisfied if the commutation relations (8.34) and (8.35) are replaced by the anticommutation relations given below.

$$
\begin{align*}
\left\{a_{k}, a_{l}^{\dagger}\right\}_{+} & =\delta_{k l} ;  \tag{8.42}\\
\left\{a_{k}, a_{l}\right\}_{+} & =0=\left\{a_{k}^{\dagger}, a_{l}^{\dagger}\right\}_{+} . \tag{8.43}
\end{align*}
$$

All the operators refer to the same time. Besides, from Eq. (8.43), we obtain

$$
\begin{equation*}
a_{k} a_{k}=0=a_{k}^{\dagger} a_{k}^{\dagger}, \tag{8.44}
\end{equation*}
$$

such that

$$
\begin{equation*}
N_{k}^{2}=a_{k}^{\dagger} a_{k} a_{k}^{\dagger} a_{k}=a_{k}^{\dagger}\left(1-a_{k}^{\dagger} a_{k}\right) a_{k}=N_{k}, \tag{8.45}
\end{equation*}
$$

[^59]which, in turn, yields the relation
\[

$$
\begin{equation*}
N_{k}\left(N_{k}-1\right)=0 . \tag{8.46}
\end{equation*}
$$

\]

This condition restricts the eigenvalue $n_{k}$ of $N_{k}$ to

$$
\begin{equation*}
n_{k}=0,1 \tag{8.47}
\end{equation*}
$$

The following results follow from the above relations.

$$
\begin{align*}
\left|n_{1}, n_{2}, \ldots, n_{k}, \ldots\right\rangle & =\left(a_{1}^{\dagger}\right)^{n_{1}}\left(a_{2}^{\dagger}\right)^{n_{2}} \ldots\left(a_{k}^{\dagger}\right)^{n_{k}} \ldots|0\rangle ;  \tag{8.48}\\
a_{k}\left|n_{1}, n_{2}, \ldots, n_{k}, \ldots\right\rangle & =(-1)^{S_{k}} n_{k}\left|n_{1}, n_{2}, \ldots, n_{k}-1, \ldots\right\rangle ;  \tag{8.49}\\
a_{k}^{\dagger}\left|n_{1}, n_{2}, \ldots, n_{k}, \ldots\right\rangle & =(-1)^{S_{k}}\left(1-n_{k}\right)\left|n_{1}, n_{2}, \ldots, n_{k}+1, \ldots\right\rangle ;(8.50) \\
N_{k}\left|n_{1}, n_{2}, \ldots, n_{k}, \ldots\right\rangle & =n_{k}\left|n_{1}, n_{2}, \ldots, n_{k}, \ldots\right\rangle ; \tag{8.51}
\end{align*}
$$

with

$$
S_{k}=\sum_{r=1}^{k-1} n_{r}
$$

Eqs. (8.49) and (8.50) state that the annihilation operator acting on an empty state or the creation operator acting on a filled state yields zero.

Since, for the Fermions, $n_{k}$ can assume only two values 0 and 1 , a $2 \times 2$ matrix representation can be given to the occupation number operator $N_{k}$.

$$
N_{k}=\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right] .
$$

Consequently, the matrix representation for the annihilation and creation operators become

$$
a_{k}=\left[\begin{array}{cc}
0 & 1 \\
0 & 0
\end{array}\right] ; \quad a_{k}^{\dagger}=\left[\begin{array}{cc}
0 & 0 \\
1 & 0
\end{array}\right]
$$

The Hamiltonian, the total energy and the total number of particles of a system of Fermions that the field represents are given by

$$
H=\sum_{k} N_{k} \epsilon_{k} ; \quad E=\sum_{k} n_{k} \epsilon_{k} ; \quad N=\sum_{k} n_{k} .
$$

Since the total number operator $N=\sum_{k} N_{k}$ commutes with the Hamiltonian,

$$
[N, H]=\sum_{k l}\left[N_{k}, N_{l}\right] \epsilon_{l}=0,
$$

it follows that the total number of Fermions in the field is conserved.
It can be shown that the anticommutation relations between the annihilation and creation operators for the Fermion field arise from the anticommutation relations between the field operator $\psi(\boldsymbol{x}, t)$ and its conjugate momentum $\Pi(\boldsymbol{x}, t)$.

$$
\begin{align*}
\left\{\psi(\boldsymbol{x}, t), \Pi\left(\boldsymbol{x}^{\prime}, t\right)\right\}_{+} & =i \hbar \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)  \tag{8.52}\\
\left\{\psi(\boldsymbol{x}, t), \psi\left(\boldsymbol{x}^{\prime}, t\right)\right\}_{+} & =0=\left\{\Pi(\boldsymbol{x}, t), \Pi\left(\boldsymbol{x}^{\prime}, t\right)\right\}_{+} \tag{8.53}
\end{align*}
$$

This anticommutation relations have no classical analogue.
We can also represent the particle number operator $N$ in terms of the field variable $\psi(\boldsymbol{x}, t)$, using the relation (8.29).

$$
\begin{align*}
N & =\sum_{k} a_{k}^{\dagger} a_{k}=\sum_{k} \int_{V} u_{k}(\boldsymbol{x}) \psi^{\dagger}(\boldsymbol{x}, t) u_{k}^{*}\left(\boldsymbol{x}^{\prime}\right) \psi\left(\boldsymbol{x}^{\prime}, t\right) d^{3} x d^{3} x^{\prime} \\
& =\int_{V} \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right) \psi^{\dagger}(\boldsymbol{x}, t) \psi\left(\boldsymbol{x}^{\prime}, t\right) d^{3} x d^{3} x^{\prime}, \quad \text { summing over } k \\
& =\int_{V} \psi^{\dagger}(\boldsymbol{x}, t) \psi(\boldsymbol{x}, t) d^{3} x \tag{8.54}
\end{align*}
$$

### 8.3 Relativistic fields

In particle mechanics, the space and time coordinates play different roles. The space coordinates are the mechanical variables whereas the time is a parameter. But in the case of classical fields, the space and time coordinates are treated similarly and they are the parameters specifying the space-time continuum, with which the field variables are defined. This facilitates the relativistically covariant formulation of the classical fields which are commonly known as relativistic fields that are manifestly Lorentz covariant.

With the introduction of the fourth coordinate $x_{0}=c t$ in the place of time coordinate, the action integral $\mathscr{A}$ defined by Eq. (8.14) changes only by a multiplicative factor but does not affect the formulation of the Hamilton's action principle which yields the same Lagrange equations of motion (8.19). Please note that the only term

$$
\frac{\partial}{\partial x_{\mu}} \frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \mathscr{L}}{\partial x_{\mu}}\right)}
$$

that depends on $x_{\mu}$ in Eq. (8.19) is not affected by the change of scale of any $x_{\mu}$.

In relativistic field theory, we should expect that all the quantities that define the field and the associated equations should be covariant under Lorentz transformation. This means that the quantities should be either world scalars or world vectors or world tensors. The action integral $\mathscr{A}$ is a world scalar. Similarly, the Lagrangian density $\mathscr{L}$ and the Hamiltonian density $\mathscr{H}$ are world scalars. The four-dimensional volume element $d^{4} x$ is invariant under Lorentz transformation. Let us show that the second rank tensor constructed from the Lagrangian density using the Lagrange equations of motion is a world tensor of second rank, the components of which yields the momentum and energy density of the field.

Given the Lagrangian density

$$
\mathscr{L}\left(\psi, \frac{\partial \psi}{\partial x_{\mu}}, x_{\mu}\right),
$$

the total derivative of the Lagrangian density is given by

$$
\begin{equation*}
\frac{d \mathscr{L}}{d x_{\mu}}=\frac{\partial \mathscr{L}}{\partial \psi} \frac{\partial \psi}{\partial x_{\mu}}+\frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \psi}{\partial x_{\nu}}\right)} \frac{\partial}{\partial x_{\mu}}\left(\frac{\partial \psi}{\partial x_{\nu}}\right)+\frac{\partial \mathscr{L}}{\partial x_{\mu}} . \tag{8.55}
\end{equation*}
$$

From Euler-Lagrange equation (8.19), we obtain

$$
\frac{\partial \mathscr{L}}{\partial \psi}=\frac{\partial}{\partial x_{\nu}} \frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \psi}{\partial x_{\nu}}\right)}
$$

. Substituting this into Eq. (8.55), we get

$$
\begin{align*}
\frac{d \mathscr{L}}{d x_{\mu}} & =\left(\frac{\partial}{\partial x_{\nu}} \frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \psi}{\partial x_{\nu}}\right)}\right) \frac{\partial \psi}{\partial x_{\mu}}+\frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \psi}{\partial x_{\nu}}\right)} \frac{\partial}{\partial x_{\mu}}\left(\frac{\partial \psi}{\partial x_{\nu}}\right)+\frac{\partial \mathscr{L}}{\partial x_{\mu}} \\
& =\frac{d}{d x_{\nu}}\left\{\frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \psi}{\partial x_{\nu}}\right)} \frac{\partial \psi}{\partial x_{\mu}}\right\}+\frac{\partial \mathscr{L}}{\partial x_{\mu}} . \tag{8.56}
\end{align*}
$$

On rearrangement, we get

$$
\begin{equation*}
\frac{d}{d x_{\nu}}\left\{\frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \psi}{\partial x_{\nu}}\right)} \frac{\partial \psi}{\partial x_{\mu}}-\mathscr{L} \delta_{\mu \nu}\right\}=-\frac{\partial \mathscr{L}}{\partial x_{\mu}} . \tag{8.57}
\end{equation*}
$$

The quantity within the curly bracket is a second rank tensor, denoted by the symbol $T_{\mu \nu}$.

$$
\begin{equation*}
T_{\mu \nu}=\frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \psi}{\partial x_{\nu}}\right)} \frac{\partial \psi}{\partial x_{\mu}}-\mathscr{L} \delta_{\mu \nu} . \tag{8.58}
\end{equation*}
$$

If the Lagrangian density $\mathscr{L}$ does not depend explicitly on $x_{\mu}$, i.e. if

$$
\frac{\partial \mathscr{L}}{\partial x_{\mu}}=0
$$

then $\mathscr{L}$ represents a free field. There are no sources or sinks or any interaction with other fields or particles. It follows from Eq. (8.57) that for a free field

$$
\begin{equation*}
\sum_{\nu} \frac{d T_{\mu \nu}}{d x_{\nu}}=0 . \tag{8.59}
\end{equation*}
$$

It can easily be verified that

$$
\begin{align*}
T_{00} & =\frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \psi}{\partial x_{0}}\right)} \frac{\partial \psi}{\partial x_{0}}-\mathscr{L} \\
& =\frac{\partial \mathscr{L}}{\partial \dot{\psi}} \dot{\psi}-\mathscr{L} . \tag{8.60}
\end{align*}
$$

This is the energy density $\mathscr{H}(\mathbf{x})$ of the free field as can be seen from Eqs. (8.21) and (8.22).

The other components of the tensor are

$$
\begin{align*}
T_{k 0} & =\frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \psi}{\partial x_{0}}\right)} \frac{\partial \psi}{\partial x_{k}}, \quad(k=1,2,3) \\
& =c \frac{\partial \mathscr{L}}{\partial \dot{\psi}} \frac{\partial \psi}{\partial x_{k}} \\
& =c \Pi \frac{\partial \psi}{\partial x_{k}}, \quad \text { using Eq. (8.21). } \tag{8.61}
\end{align*}
$$

Defining the momentum density $\mathscr{P}$ by

$$
\begin{equation*}
\mathscr{P}_{k}=\Pi \frac{\partial \psi}{\partial x_{k}}, \tag{8.62}
\end{equation*}
$$

we find that

$$
\begin{equation*}
T_{k 0}=c \mathscr{P}_{k} . \tag{8.63}
\end{equation*}
$$

## Review Questions

8.1 Discuss briefly the Lagrangian-Hamiltonian formalism in classical mechanics and deduce the Euler-Lagrange equations of motion and also Hamilton's canonical equations of motion. Explain how these concepts can be applied to the classical fields.
8.2 What is meant by the second quantization? Discuss how the Schrödinger field can be quantized to yield a system of Bosons. Define the number operator and express the Hamiltonian in terms of the number operator.
8.3 Discuss how the Schrödinger equation can be quantized to yield a system of Fermions. Show that the number operator can assume only two values, 0 or 1. Obtain an expression for the Hamiltonian in terms of the number operator.
8.4 What are the features that a relativistic field should exhibit? Given the Lagrangian density, construct a second rank tensor using Euler-Lagrange equations of motion and show that the components of the second rank tensor yield the momentum and energy density of the field.

## Problems

8.1 Show that in the case of motion of a single particle, the principle of least action yields the second law of Newton.
8.2 Show that the Lagrangian density for the Schrödinger field which obeys Schrödinger equation is

$$
\mathscr{L}_{S}=i \hbar \phi^{*} \dot{\phi}-\frac{\hbar^{2}}{2 m} \boldsymbol{\nabla} \phi^{*} \cdot \boldsymbol{\nabla} \phi-V \phi^{*} \phi .
$$

8.3 Given the Lagrangian density

$$
\mathscr{L}=\frac{1}{2}\left\{\dot{\phi}^{2}-(\boldsymbol{\nabla} \phi)^{2}-m^{2} \phi^{2}\right\},
$$

deduce the equation for the field $\phi$ that the Lagrangian density represents.

## Solutions to Problems

8.1 Newton's second law states that a point particle of mass $m$ moving in a potential field $V(x)$ experiences a force $F$ resulting in an acceleration $\frac{d^{2} x}{d t^{2}}$ of the particle.

$$
m \frac{d^{2} x}{d t^{2}}=F=-\frac{d V(x)}{d t} .
$$

The Lagrangian $L$ of the particle is given by

$$
L=T-V=\frac{1}{2} m\left(\frac{d x}{d t}\right)^{2}-V(x)
$$

where $T$ is the kinetic energy. The action $\mathscr{A}$ is defined by

$$
\mathscr{A}=\int_{t_{1}}^{t_{2}} L(x, \dot{x}) d t,
$$



Figure 8.1: Multiple paths available for the particle between two space-time points A and B , of which the particle chooses one for which the action $\mathscr{A}$ is minimal.
where the integral is taken over the entire path of the particle from time $t_{1}$ to $t_{2}$ as shown in Fig. 8.1. There is an infinite number of possible paths for the particle but what the actual path that the particle will take is given by the principle of least action. In other words, the particle chooses a path for which the action $\mathscr{A}$ is minimum.
Consider a variation in the path

$$
x(t) \rightarrow x^{\prime}(t)=x(t)+a(t), \quad a \ll x .
$$

Since the end points of the path are fixed,

$$
a\left(t_{1}\right)=a\left(t_{2}\right)=0
$$

On the substitution $x \rightarrow x^{\prime}$, the action $\mathscr{A}$ becomes

$$
\begin{aligned}
\mathscr{A} \rightarrow \mathscr{A}^{\prime} & =\int_{t_{1}}^{t_{2}}\left[\frac{1}{2} m(\dot{x}+\dot{a})^{2}-V(x+a)\right] d t \\
& =\int_{t_{1}}^{t_{2}}\left[\frac{1}{2} m \dot{x}^{2}+m \dot{x} \dot{a}-\left\{V(x)+a V^{\prime}(x)\right\}\right] d t+O\left(a^{2}\right) \\
& =\mathscr{A}+\int_{t_{1}}^{t_{2}}\left[m \dot{x} \dot{a}-a V^{\prime}(x)\right] d t \\
& =\mathscr{A}+\delta \mathscr{A}
\end{aligned}
$$

where

$$
\delta \mathscr{A}=\int_{t_{1}}^{t_{2}}\left[m \dot{x} \dot{a}-a V^{\prime}(x)\right] d t, \quad V^{\prime}(x)=\frac{d V(x)}{d x}
$$

If $\mathscr{A}$ is a minimum under the variation in $x$, then $\delta \mathscr{A}=0$. Integrating the first term in $\delta \mathscr{A}$ by parts, we get

$$
\int_{t_{1}}^{t_{2}} \dot{x} \dot{a} d t=\left.\dot{x} a\right|_{t_{1}} ^{t_{2}}-\int_{t_{1}}^{t_{2}} a \ddot{x} d t=-\int_{t_{1}}^{t_{2}} a \ddot{x} d t
$$

since $a\left(t_{1}\right)=a\left(t_{2}\right)=0$ because of the fixed end points. Thus we get

$$
\delta \mathscr{A}=-\int_{t_{1}}^{t_{2}}\left[m a \ddot{x}+a V^{\prime}(x)\right] d t=0
$$

which is satisfied if

$$
m \ddot{x}=-V^{\prime}(x)
$$

Thus we have deduced Newton's second law of motion from the principle of least action.
8.2 Given the Lagrangian density $\mathscr{L}_{S}$ of the field, find the field equation which satisfies the Euler-Lagrange equation (8.19).
Let us write explicitly the Euler-Lagrange equation.

$$
\begin{align*}
\frac{\partial \mathscr{L}}{\partial \phi}-\sum_{\mu=0}^{3} \frac{\partial}{\partial x^{\mu}} \frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \phi}{\partial x^{\mu}}\right)} & =0 \\
\text { or } \quad \frac{\partial \mathscr{L}}{\partial \phi}-\frac{\partial}{\partial t} \frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \phi}{\partial t}\right)}-\sum_{k=1}^{3} \frac{\partial}{\partial x_{k}} \frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \phi}{\partial x^{k}}\right)} & =0 \tag{8.64}
\end{align*}
$$

Given the Lagrangian density

$$
\mathscr{L}_{s}=i \hbar \phi^{*} \dot{\phi}-\frac{\hbar^{2}}{2 m} \nabla \phi^{*} \cdot \nabla \phi-V \phi^{8} \phi
$$

we find

$$
\begin{align*}
\frac{\partial \mathscr{L}_{S}}{\partial \phi} & =-V \phi^{*} & & \\
\frac{\partial \mathscr{L}_{S}}{\partial \boldsymbol{\nabla} \phi} & =-\frac{\hbar^{2}}{2 m} \nabla \phi^{*}, & & \nabla \cdot\left(\frac{\partial \mathscr{L}_{S}}{\partial \boldsymbol{\nabla} \phi}\right)=-\frac{\hbar^{2}}{2 m} \nabla^{2} \phi^{*}  \tag{8.65}\\
\frac{\partial \mathscr{L}_{S}}{\partial \dot{\phi}} & =i \hbar \phi^{*}, & & \frac{\partial}{\partial t}\left(\frac{\partial \mathscr{L}_{S}}{\partial \dot{\phi}}\right)=i \hbar \frac{\partial \phi^{*}}{\partial t}
\end{align*}
$$

Substituting (8.65) into (8.64), we get

$$
-\frac{\hbar^{2}}{2 m} \nabla^{2} \phi^{*}+V \phi^{*}+i \hbar \frac{\partial \phi^{*}}{\partial t}=0
$$

the complex conjugate of which is the Schrödinger equation.

$$
-\frac{\hbar^{2}}{2 m} \nabla^{2} \phi+V \phi-i \hbar \frac{\partial \phi}{\partial t}=0
$$

8.3 Given the Lagrangian density

$$
\mathscr{L}=\frac{1}{2}\left\{\dot{\phi}^{2}-(\boldsymbol{\nabla} \phi)^{2}-m^{2} \phi^{2}\right\}
$$

it can be written in the four-component space-time notation.

$$
\mathscr{L}=\frac{1}{2} \sum_{\mu}\left(\frac{\partial \phi}{\partial x^{\mu}}\right)^{2}-\frac{1}{2} m^{2} \phi^{2} .
$$

Since the Lagrangian density is given, we need to find the field equation that satisfies the Euler-Lagrangian equation.

$$
\frac{\partial \mathscr{L}}{\partial \phi}-\sum_{\mu} \frac{\partial}{\partial x^{\mu}} \frac{\partial \mathscr{L}}{\partial \frac{\partial \phi}{\partial x^{\mu}}}=0 .
$$

Since

$$
\begin{aligned}
\frac{\partial \mathscr{L}}{\partial \phi} & =-m^{2} \phi \\
\frac{\partial \mathscr{L}}{\partial \frac{\partial \phi}{\partial x^{\mu}}} & =\frac{\partial \phi}{\partial x^{\mu}} \\
\frac{\partial}{\partial x^{\mu}} \frac{\partial \mathscr{L}}{\partial \frac{\partial \phi}{\partial x^{\mu}}} & =\frac{\partial^{2} \phi}{\partial x^{\mu^{2}}}
\end{aligned}
$$

Substituting these values in Euler-lagrange equation, we get the field equation which can be identified as the Klein-Gordon equation

$$
\left(\square+m^{2}\right) \phi(\mathbf{x})=0,
$$

where

$$
\square=\frac{\partial}{\partial x^{\mu}} \frac{\partial}{\partial x_{\mu}}=\frac{\partial^{2}}{\partial t^{2}}-\nabla^{2},
$$

in natural units.

## Chapter 9

## The Scalar Fields

The well-known equation for the scalar field is the Klein-Gordon (K-G) equation. The Klein-Gordon equation is a relativistic wave equation, obtained earlier for a single particle but it was initially unacceptable for the following two reasons. It admitted not only negative energy states but also negative probability densities. The latter was a more serious defect and was physically unacceptable. Only when Pauli and Weisskopf ${ }^{1}$ found a way out and reinterpreted the Klein-Gordon equation as a field equation in the same sense as Maxwell's equations for electromagnetic field, the interest revived and became a subject of intense study.

Klein-Gordon equations are used to describe particles of spin 0 such as $\pi$ mesons and $K$ mesons. Real i.e. Hermitian fields represent uncharged particles and complex (non-Hermitian) fields represent charged particles.

### 9.1 One-component real field

First let us write down the K-G equation in different forms as are found in the literature.

$$
\begin{equation*}
\left(\square+\kappa^{2}\right) \Phi(\mathbf{x})=0 ; \quad\left(\partial^{\mu} \partial_{\mu}+\kappa^{2}\right) \Phi(\mathbf{x})=0 ; \tag{9.1}
\end{equation*}
$$

where

$$
\begin{equation*}
\square=\frac{1}{c^{2}} \frac{\partial^{2}}{\partial t^{2}}-\nabla^{2} . \tag{9.2}
\end{equation*}
$$

[^60]The symbols $\partial^{\mu}, \partial_{\mu}$ and $\kappa$ denote

$$
\begin{align*}
\partial^{\mu} & =\frac{\partial}{\partial x_{\mu}} ; & & \partial_{\mu}=\frac{\partial}{\partial x^{\mu}} ; \quad(\mu=0,1,2,3)  \tag{9.3}\\
\kappa & =\frac{m c}{\hbar} . & & \text { (Inverse Compton wavelength of the particle) } \tag{9.4}
\end{align*}
$$

Hereafter, we shall use natural units $(\hbar=c=1)$. In natural units, the K-G equation (9.1) can be rewritten as

$$
\begin{equation*}
\left(\partial^{\mu} \partial_{\mu}+m^{2}\right) \Phi(\mathbf{x})=0 \tag{9.5}
\end{equation*}
$$

The plane wave solution of K-G Eq. $\left(\partial^{\mu} \partial_{\mu}+m^{2}\right) u_{k}(\mathbf{x})=0$ is

$$
\begin{equation*}
u_{k}(\mathbf{x})=\frac{1}{\sqrt{V}} \frac{1}{\sqrt{2 \omega_{k}}} e^{-i \mathbf{k} \cdot \mathbf{x}} \tag{9.6}
\end{equation*}
$$

where $V$ is the quantization volume and

$$
\begin{align*}
\mathbf{k} \cdot \mathbf{x} & =k_{0} x_{0}-\boldsymbol{k} \cdot \boldsymbol{x}=\omega_{k} t-\boldsymbol{k} \cdot \boldsymbol{x}  \tag{9.7}\\
\mathbf{k}^{2} & =k^{\mu} k_{\mu}=\omega_{k}^{2}-k^{2}=m^{2}  \tag{9.8}\\
\omega_{k}^{2} & =k^{2}+m^{2} \tag{9.9}
\end{align*}
$$

We use the convention of representing the four-vectors by bold up-right letters ( $\mathbf{k}, \mathbf{x}$ ), three vectors by bold italics $(\boldsymbol{k}, \boldsymbol{x})$ and the scalars by ordinary italics $(k, x)$.

### 9.1.1 Fourier decomposition of the field

The field operator $\Phi(\mathbf{x})$ in Eq. (9.5) is expanded in terms of the complete set of $u_{k}(\mathbf{x})$, defined by Eq. (9.6).

$$
\begin{align*}
\Phi(\mathbf{x}) & =\sum_{k}\left(a_{k} u_{k}(\mathbf{x})+a_{k}^{\dagger} u_{k}^{*}(\mathbf{x})\right) \\
& =\frac{1}{\sqrt{V}} \sum_{k} \frac{1}{\sqrt{2 \omega_{k}}}\left(a(\boldsymbol{k}) e^{-i \mathbf{k} \cdot \mathbf{x}}+a^{\dagger}(\boldsymbol{k}) e^{i \mathbf{k} \cdot \mathbf{x}}\right) \tag{9.10}
\end{align*}
$$

The second term on the right hand side of Eq. (9.10) is the hermitian conjugate of the first term and is included in order to make the field function $\Phi(\mathbf{x})$ real. The expansion coefficients ${ }^{2} a_{k}$ and $a_{k}^{\dagger}$ will have to be interpreted as the annihilation and creation operators and the field

[^61]operator $\Phi(\mathbf{x})$ represents an assembly of particles. More often, Eq. (9.10) is split into two terms
\[

$$
\begin{equation*}
\Phi(\mathbf{x})=\Phi^{+}(\mathbf{x})+\Phi^{-}(\mathbf{x}) \tag{9.11}
\end{equation*}
$$

\]

where $\Phi^{+}(\mathbf{x})$ which carries the annihilation operator $a_{k}$ corresponds to the positive energy part of the decomposition of the field function $\Phi(\mathbf{x})$ and $\Phi^{-}(\mathbf{x})$ which carries the creation operator $a_{k}^{\dagger}$ corresponds to the negative energy part of the decomposition of the field function $\Phi(\mathbf{x})$. Using the prescription

$$
\frac{1}{\sqrt{V}} \sum_{k} \longrightarrow \frac{1}{(2 \pi)^{3 / 2}} \int d^{3} k
$$

we can make a transition from the discrete values of $\boldsymbol{k}$ to continuous values and write down explicitly $\Phi^{+}(\mathbf{x})$ and $\Phi^{-}(\mathbf{x})$.

$$
\begin{align*}
& \Phi^{+}(\mathbf{x})=\frac{1}{(2 \pi)^{3 / 2}} \int_{k_{0}>0} \frac{d^{3} k}{\sqrt{2 \omega_{k}}} a(\boldsymbol{k}) e^{-i \mathbf{k} \cdot \mathbf{x}}  \tag{9.12}\\
& \Phi^{-}(\mathbf{x})=\frac{1}{(2 \pi)^{3 / 2}} \int_{k_{0}>0} \frac{d^{3} k}{\sqrt{2 \omega_{k}}} a^{\dagger}(\boldsymbol{k}) e^{i \mathbf{k} \cdot \mathbf{x}} \tag{9.13}
\end{align*}
$$

The splitting of the field function $\Phi$ into two parts $\Phi^{+}(\mathbf{x})$ and $\Phi^{-}(\mathbf{x})$ that correspond to positive and negative frequency parts, is introduced here, since we will have many occasions to use them at a later stage in our discussion.

Since we are using the Lagrangian formalism, the Lagrangian density will play a central role and so let us conjecture the Lagrangian density $\mathscr{L}_{K G}$ that will satisfy the Euler-Lagrange equation and yield the KleinGordon equation.

$$
\begin{align*}
\mathscr{L}_{K G}(\mathbf{x}) & =\frac{1}{2}\left\{\sum_{\mu}\left(\frac{\partial \Phi}{\partial x_{\mu}}\right)^{2}-m^{2} \Phi^{2}\right\} \\
& =\frac{1}{2}\left(\partial^{\mu} \Phi(\mathbf{x}) \partial_{\mu} \Phi(\mathbf{x})-m^{2} \Phi^{2}(\mathbf{x})\right) \tag{9.14}
\end{align*}
$$

The repeated index $\mu$ implies summation over $\mu$. Given the above Lagrangian density, one can find the conjugate momentum density $\Pi(\mathbf{x})$ and the Hamiltonian density $\mathscr{H}(\mathbf{x})$ by the established procedure.

$$
\begin{equation*}
\Pi(\mathbf{x})=\frac{\partial \mathscr{L}_{K G}}{\partial \dot{\Phi}}=\partial_{0} \Phi=\dot{\Phi} \tag{9.15}
\end{equation*}
$$

$$
\begin{align*}
\mathscr{H}(\mathbf{x}) & =\Pi(\mathbf{x}) \dot{\Phi}(\mathbf{x})-\mathscr{L}_{K G}(\mathbf{x}) \\
& =\partial^{0} \Phi \partial_{0} \Phi-\frac{1}{2}\left(\partial^{\mu} \Phi(\mathbf{x}) \partial_{\mu} \Phi(\mathbf{x})-m^{2} \Phi^{2}(\mathbf{x})\right) \\
& =\frac{1}{2}\left\{\dot{\Phi}^{2}+(\boldsymbol{\nabla} \Phi)^{2}+m^{2} \Phi^{2}\right\} . \tag{9.16}
\end{align*}
$$

In Eq. (9.16), $\boldsymbol{\nabla}$ denotes the ordinary three-dimensional gradient. The energy density as given by Eq. (9.16) is a positive-definite quantity as should be expected.

The total energy of the field is ${ }^{3}$

$$
\begin{align*}
H & =\int \mathscr{H} d^{3} x=\int\left(\Pi \dot{\Phi}-\mathscr{L}_{K G}\right) d^{3} x \\
& =\int\left(\frac{\partial \mathscr{L}}{\partial \dot{\Phi}} \dot{\Phi}-\mathscr{L}\right) d^{3} x=\int\left(\frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \Phi}{\partial x_{0}}\right)} \frac{\partial \Phi}{\partial x_{0}}-\mathscr{L}\right) d^{3} x \\
& =\int T_{00} d^{3} x \tag{9.17}
\end{align*}
$$

where $T_{00}$ is the $(0,0)$ component of the canonical energy momentum tensor $T_{\mu \nu}$ defined earlier in Eq. (8.58).

$$
\begin{equation*}
T_{\mu \nu}=\frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \Phi}{\partial x_{\nu}}\right)} \frac{\partial \Phi}{\partial x_{\mu}}-\mathscr{L} \delta_{\mu \nu} . \tag{9.18}
\end{equation*}
$$

Using the Fourier decomposition (9.10) of the field, the conjugate momentum density $\Pi(\mathbf{x})$ given by Eq. (9.15) can be written in the expanded form using annihilation and creation operators $a(\boldsymbol{k})$ and $a^{\dagger}(\boldsymbol{k})$.

$$
\begin{equation*}
\Pi(\mathbf{x})=\dot{\Phi}(\mathbf{x})=\frac{1}{\sqrt{V}} \sum_{k} \frac{-i \omega_{k}}{\sqrt{2 \omega_{k}}}\left(a(\boldsymbol{k}) e^{-i \mathbf{k} \cdot \mathbf{x}}-a^{\dagger}(\boldsymbol{k}) e^{i \mathbf{k} \cdot \mathbf{x}}\right) . \tag{9.19}
\end{equation*}
$$

Since the commutation relations between $\Phi$ and $\Pi$ and the Hamiltonian $H$ do not depend on time, let us choose $t=0$. Then,

$$
\begin{equation*}
\Phi(\boldsymbol{x}, 0)=\frac{1}{\sqrt{V}} \sum_{k} \frac{1}{\sqrt{2 \omega_{k}}}\left(a(\boldsymbol{k}) e^{i \boldsymbol{k} \cdot \boldsymbol{x}}+a^{\dagger}(\boldsymbol{k}) e^{-i \boldsymbol{k} \cdot \boldsymbol{x}}\right) . \tag{9.20}
\end{equation*}
$$

[^62]A transition from the discrete values of $k$ to continuous values of $k$ can be made by using the prescription (vide Solved Problem (9.4))

$$
\begin{gather*}
\frac{1}{\sqrt{V}} \sum_{k} \longrightarrow \frac{1}{(2 \pi)^{3 / 2}} \int d^{3} k . \\
\Phi(\boldsymbol{x}, 0)=\frac{1}{(2 \pi)^{3 / 2}} \int \frac{d^{3} k}{\sqrt{2 \omega_{k}}}\left(a(\boldsymbol{k}) e^{i \boldsymbol{k} \cdot \boldsymbol{x}}+a^{\dagger}(\boldsymbol{k}) e^{-i \boldsymbol{k} \cdot \boldsymbol{x}}\right) .  \tag{9.21}\\
\Pi(\boldsymbol{x}, 0)=\frac{1}{(2 \pi)^{3 / 2}} \int \frac{d^{3} k}{\sqrt{2 \omega_{k}}}\left(-i \omega_{k}\right)\left(a(\boldsymbol{k}) e^{i \boldsymbol{k} \cdot \boldsymbol{x}}-a^{\dagger}(\boldsymbol{k}) e^{-i \boldsymbol{k} \cdot \boldsymbol{x}}\right) . \tag{9.22}
\end{gather*}
$$

Taking the Fourier inversions of (9.21) and (9.22), we obtain (vide Solved Problem (9.2))

$$
\begin{align*}
a(\boldsymbol{k}) & =\frac{1}{(2 \pi)^{3 / 2}} \int \frac{d^{3} x}{\sqrt{2 \omega_{k}}} e^{-i \boldsymbol{k} \cdot \boldsymbol{x}}\left\{\omega_{k} \Phi(\boldsymbol{x})+i \Pi(\boldsymbol{x})\right\}  \tag{9.23}\\
a^{\dagger}(\boldsymbol{k}) & =\frac{1}{(2 \pi)^{3 / 2}} \int \frac{d^{3} x}{\sqrt{2 \omega_{k}}} e^{i \boldsymbol{k} \cdot \boldsymbol{x}}\left\{\omega_{k} \Phi(\boldsymbol{x})-i \Pi(\boldsymbol{x})\right\} \tag{9.24}
\end{align*}
$$

### 9.1.2 Quantization of the scalar field

The scalar field is quantized by invoking the equal-time commutation relations between the field function $\Phi(\mathbf{x})$ and its conjugate momentum function $\Pi(\mathrm{x})$.

$$
\begin{align*}
& {\left[\Phi(\boldsymbol{x}, t), \Pi\left(\boldsymbol{x}^{\prime}, t\right)\right]_{-}=i \hbar \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)=i \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)}  \tag{9.25}\\
& {\left[\Phi(\boldsymbol{x}, t), \Phi\left(\boldsymbol{x}^{\prime}, t\right)\right]_{-}=\left[\Pi(\boldsymbol{x}, t), \Pi\left(\boldsymbol{x}^{\prime}, t\right)\right]_{-}=0} \tag{9.26}
\end{align*}
$$

They, in turn, yield the commutation relations between the annihilation and creation operators. (vide Solved Problem (9.3))

$$
\begin{align*}
{\left[a(\boldsymbol{k}), a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-} } & =\delta_{\boldsymbol{k}, \boldsymbol{k}^{\prime}} ;  \tag{9.27}\\
{\left[a(\boldsymbol{k}), a\left(\boldsymbol{k}^{\prime}\right)\right]_{-} } & =\left[a^{\dagger}(\boldsymbol{k}), a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=0 . \tag{9.28}
\end{align*}
$$

There is a very close analogy between the simple harmonic oscillator ${ }^{4}$ and the field. The operators $a(\boldsymbol{k})$ and $a^{\dagger}(\boldsymbol{k})$ are interpreted as the annihilation and creation operators for a field quantum with momentum $\boldsymbol{k}$. The number operator $N(\boldsymbol{k})$ is defined as

$$
\begin{equation*}
N(\boldsymbol{k})=a^{\dagger}(\boldsymbol{k}) a(\boldsymbol{k}), \tag{9.29}
\end{equation*}
$$

[^63]which operating on the basic state in the Fock space
\[

$$
\begin{equation*}
\left|n_{1}, n_{2}, n_{3}, \cdots, n_{k}, \cdots\right\rangle \tag{9.30}
\end{equation*}
$$

\]

yields the number of quanta associated with the momentum $\boldsymbol{k}$.

$$
\begin{equation*}
N(\boldsymbol{k})\left|n_{1}, n_{2}, n_{3}, \cdots, n_{k}, \cdots\right\rangle=n_{k}\left|n_{1}, n_{2}, n_{3}, \cdots, n_{k}, \cdots\right\rangle \tag{9.31}
\end{equation*}
$$

Using Eq. (9.16), it can be shown that the total Hamiltonian $H$ is (vide Solved Problem (9.8))

$$
\begin{align*}
H & =\int \mathscr{H} d^{3} x=\frac{1}{2} \int\left(\Pi^{2}+(\boldsymbol{\nabla} \Phi)^{2}+m^{2} \Phi^{2}\right) d^{3} x \\
& =\frac{1}{2} \sum_{k} \omega_{k}\left(a^{\dagger}(\boldsymbol{k}) a(\boldsymbol{k})+a(\boldsymbol{k}) a^{\dagger}(\boldsymbol{k})\right) \\
& =\sum_{k} \omega_{k}\left(a^{\dagger}(\boldsymbol{k}) a(\boldsymbol{k})+\frac{1}{2}\right) . \tag{9.32}
\end{align*}
$$

In the last step, the commutation relation (9.27) is used. The ground state is one with $n_{k}=0$ for all $k$ and we should expect the ground state energy to be zero. But on the other hand, it becomes infinity $E=\sum_{k} \frac{1}{2} \omega_{k}$. The second quantized Klein-Gordon field thus gives a description of an assembly of infinitely many identical non-interacting particles of the same mass. They are not only spinless but also uncharged. That the associated quanta are uncharged follows from the fact that the four-vector current density corresponding to the Hermitian field is zero.

### 9.1.3 Ground state and normal ordering

The ground state of the field which is called the vacuum is defined as

$$
\begin{equation*}
a(\boldsymbol{k})|0\rangle=0, \quad \text { for all } \boldsymbol{k} . \tag{9.33}
\end{equation*}
$$

Assuming the normalization of the ground state

$$
\begin{equation*}
\langle 0 \mid 0\rangle=1, \tag{9.34}
\end{equation*}
$$

we find the expectation value of the total Hamiltonian in the ground state to be

$$
\begin{equation*}
\langle 0| H|0\rangle=\frac{1}{2} \sum_{k=0}^{\infty} \omega_{k}, \tag{9.35}
\end{equation*}
$$

which is infinite. This is because, we have one oscillator for each value of momentum and we have an infinite number of such oscillators, each having a zero point energy. Thus the ground state energy becomes infinite. Now, let us redefine the zero of energy in such a way as to make the ground state energy zero. This can be done by a simple prescription called normal ordering. Whenever, we have a product of annihilation and creation operators, simply arrange all the annihilation operators to the right of all creation operators as if the commutators were zero. This is called normal ordering. Once we have done this arrangement, then the operators are treated once again as operators with the usual commutation relations. Using this prescription of normal ordering, we find the expectation value of the normal-ordered Hamiltonian denoted by the symbol : $H$ : or $N()$ to be

$$
\begin{equation*}
\left.\langle\Psi|: H:|\Psi\rangle=\sum_{k} \omega_{k}\langle\Psi| a^{\dagger}(\boldsymbol{k}) a(\boldsymbol{k})|\Psi\rangle=\sum_{k} \omega_{k}|a(\boldsymbol{k})| \Psi\right\rangle\left.\right|^{2} . \tag{9.36}
\end{equation*}
$$

The energy thus obtained is always positive and the expectation value of the energy of the vacuum now becomes

$$
\begin{equation*}
\langle 0|: H:|0\rangle=0, \tag{9.37}
\end{equation*}
$$

so that it corresponds to the ground state energy which is zero.

### 9.2 Complex scalar field

Any complex scalar field can be represented by two independent fields, one corresponding to the real part and the other corresponding to the imaginary part or alternatively by the complex field $\Phi(\mathbf{x})$ and its complex conjugate $\Phi^{*}(\mathbf{x})$, each of which obey the Klein-Gordon equation as given below:

$$
\begin{equation*}
\left(\partial^{\mu} \partial_{\mu}+m^{2}\right) \Phi(\mathbf{x})=0 ; \quad\left(\partial^{\mu} \partial_{\mu}+m^{2}\right) \Phi^{*}(\mathbf{x})=0 \tag{9.38}
\end{equation*}
$$

The Lagrangian density for such a complex field is given by

$$
\begin{equation*}
\mathscr{L}(\mathbf{x})=\partial^{\mu} \Phi(\mathbf{x}) \partial_{\mu} \Phi^{*}(\mathbf{x})-m^{2} \Phi(\mathbf{x}) \Phi^{*}(\mathbf{x}) \tag{9.39}
\end{equation*}
$$

where $\Phi(\mathbf{x})$ and $\Phi^{*}(\mathbf{x})$ are to be treated as independent functions which can be treated as a linear combination of two independent real fields $\Phi_{1}(\mathbf{x})$ and $\Phi_{2}(\mathbf{x})$.

$$
\begin{equation*}
\Phi(\mathbf{x})=\frac{1}{\sqrt{2}}\left(\Phi_{1}(\mathbf{x})+i \Phi_{2}(\mathbf{x})\right) ; \quad \Phi^{*}(\mathbf{x})=\frac{1}{\sqrt{2}}\left(\Phi_{1}(\mathbf{x})-i \Phi_{2}(\mathbf{x})\right) . \tag{9.40}
\end{equation*}
$$

Each of these fields obey the Klein-Gordon equation.

$$
\begin{equation*}
\left(\partial^{\mu} \partial_{\mu}+m^{2}\right) \Phi_{r}(\mathbf{x})=0, \quad r=1,2 . \tag{9.41}
\end{equation*}
$$

The Lagrangian density for the complex field can also be written as the sum of the Lagrangian densities of the two independent real fields $\Phi_{1}$ and $\Phi_{2}$.

$$
\begin{equation*}
\mathscr{L}=\frac{1}{2} \sum_{r=1}^{2}\left(\partial^{\mu} \Phi_{r} \partial_{\mu} \Phi_{r}-m^{2} \Phi_{r}\right), \quad r=1,2 . \tag{9.42}
\end{equation*}
$$

The conjugate momentum density of each field is given by

$$
\begin{equation*}
\Pi_{r}(\mathbf{x})=\frac{\partial \mathscr{L}}{\partial \dot{\Phi}_{r}}=\partial_{0} \Phi_{r}=\dot{\Phi}_{r}, \quad r=1,2 \tag{9.43}
\end{equation*}
$$

The Hamiltonian density of the complex field is given by

$$
\begin{equation*}
\mathscr{H}=\sum_{r=1}^{2} \Pi_{r}(\mathbf{x}) \dot{\Phi}_{r}(\mathbf{x})-\mathscr{L}(\mathbf{x}) \tag{9.44}
\end{equation*}
$$

from which the total Hamiltonian $H$ is obtained by performing the threedimensional space integration.

$$
\begin{align*}
H & =\int \mathscr{H} d^{3} x \\
& =\sum_{r=1}^{2}\left\{\int \partial^{0} \Phi_{r} \partial_{0} \Phi_{r} d^{3} x-\frac{1}{2} \int\left(\partial^{\mu} \Phi_{r} \partial_{\mu} \Phi_{r}-m^{2}\right) d^{3} x\right\} \tag{9.45}
\end{align*}
$$

As was done earlier for the real scalar field, we can do the Fourier decomposition of the field for $\Phi_{1}$ and $\Phi_{2}$.

$$
\begin{equation*}
\Phi_{r}(\mathbf{x})=\frac{1}{\sqrt{V}} \sum_{k} \frac{1}{\sqrt{2 \omega_{k}}}\left\{a_{r}(\boldsymbol{k}) e^{-i \mathbf{k} \cdot \mathbf{x}}+a_{r}^{\dagger}(\boldsymbol{k}) e^{i \mathbf{k} \cdot \mathbf{x}}\right\} \tag{9.46}
\end{equation*}
$$

where the operators $a_{r}(\boldsymbol{k})$ and $a_{r}^{\dagger}(\boldsymbol{k})$ are the annihilation and creation operators which obey the commutation relations

$$
\begin{align*}
& {\left[a_{r}(\boldsymbol{k}), a_{s}^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=\delta_{r s} \delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right)}  \tag{9.47}\\
& {\left[a_{r}(\boldsymbol{k}), a_{s}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=0=\left[a_{r}^{\dagger}(\boldsymbol{k}), a_{s}^{\dagger}(\boldsymbol{k})\right]_{-} .} \tag{9.48}
\end{align*}
$$

The commutators indicate that the field represents a system of bosons. The complex scalar field represents two types of particles and this would become clear by looking at the Fourier expansion of $\Phi$ using the relations (9.40).

$$
\begin{equation*}
\Phi(\mathbf{x})=\frac{1}{\sqrt{V}} \sum_{k} \frac{1}{\sqrt{2 \omega_{k}}}\left\{a(\boldsymbol{k}) e^{-i \mathbf{k} \cdot \mathbf{x}}+b^{\dagger}(\boldsymbol{k}) e^{i \mathbf{k} \cdot \mathbf{x}}\right\} \tag{9.49}
\end{equation*}
$$

where

$$
\begin{equation*}
a(\boldsymbol{k})=\frac{1}{\sqrt{2}}\left(a_{1}(\boldsymbol{k})+i a_{2}(\boldsymbol{k})\right), \quad b(\boldsymbol{k})=\frac{1}{\sqrt{2}}\left(a_{1}(\boldsymbol{k})-i a_{2}(\boldsymbol{k})\right) . \tag{9.50}
\end{equation*}
$$

The operators $a(\boldsymbol{k})$ and $b(\boldsymbol{k})$ and their Hermitian conjugates satisfy the commutation relations

$$
\begin{equation*}
\left[a(\boldsymbol{k}), a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=\left[b(\boldsymbol{k}), b^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=\delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right) \tag{9.51}
\end{equation*}
$$

All other commutators with different combinations vanish.
From Eq. (9.49), we can at once write down the Fourier expansion for $\Phi^{*}(\mathrm{x})$.

$$
\begin{equation*}
\Phi^{*}(\mathbf{x})=\frac{1}{\sqrt{V}} \sum_{k} \frac{1}{\sqrt{2 \omega_{k}}}\left\{a^{\dagger}(\boldsymbol{k}) e^{i \mathbf{k} \cdot \mathbf{x}}+b(\boldsymbol{k}) e^{-i \mathbf{k} \cdot \mathbf{x}}\right\} \tag{9.52}
\end{equation*}
$$

A complex scalar field is, in fact, consists of two fields $\Phi_{1}(\mathbf{x})$ and $\Phi_{2}(\mathbf{x})$ and $a_{1}^{\dagger}$ and $a_{2}^{\dagger}$ are the creation operators that produce the quanta that correspond to these fields. It is clear that a complex scalar field represents two different sets of particles. Instead of the fields $\Phi_{1}(\mathbf{x})$ and $\Phi_{2}(\mathbf{x})$, we can as well choose $\Phi(\mathbf{x})$ and $\Phi^{*}(\mathbf{x})$. For the latter choice, $a^{\dagger}$ and $b^{\dagger}$ are the creation operators for the field quanta that are described by the fields $\Phi(\mathbf{x})$ and $\Phi^{*}(\mathbf{x})$. The latter choice is preferable since the Lagrangian (9.39) is invariant under the transformation

$$
\begin{equation*}
\Phi \rightarrow \Phi e^{i e \alpha}, \quad \Phi^{*} \rightarrow \Phi^{*} e^{-i e \alpha} \tag{9.53}
\end{equation*}
$$

This transformation leads to the conservation of charge and current of the fields.

### 9.2.1 Charge-current density

Define the four vector current $j^{\mu}(\mathbf{x})$ by

$$
\begin{align*}
j^{\mu}(\mathbf{x}) & =-i e\left(\frac{\partial \mathscr{L}}{\partial \Phi_{, \mu}} \Phi-\frac{\partial \mathscr{L}}{\partial \Phi_{, \mu}^{*}} \Phi^{*}\right) \\
& =i e\left\{\left(\partial_{\mu} \Phi^{*}\right) \Phi-\left(\partial_{\mu} \Phi\right) \Phi^{*}\right\} \tag{9.54}
\end{align*}
$$

Using the Klein-Gordon field Eqs. (9.38), we get

$$
\begin{equation*}
\partial_{\mu} j^{\mu}(\mathbf{x})=i e m^{2}\left[\Phi^{*}(\mathbf{x}), \Phi(\mathbf{x})\right]_{-}=0 \tag{9.55}
\end{equation*}
$$

The four-vector $j^{\mu}(\mathbf{x})$ is to be interpreted as the charge-current density four vector and $j^{0}$ is the electric charge density.

The total charge of the field is

$$
\begin{align*}
Q & =\int j^{0}(\mathbf{x}) d^{3} x \\
& =e \int\left\{\left(\partial_{0} \Phi^{*}\right) \Phi-\left(\partial_{0} \Phi\right) \Phi^{*}\right\} d^{3} x, \quad \text { using Eq. }(9.54) \\
& =\frac{e}{2} \sum_{k k^{\prime}}\left[\left\{a\left(\boldsymbol{k}^{\prime}\right), a^{\dagger}(\boldsymbol{k})\right\}_{+}-\left\{b^{\dagger}(\boldsymbol{k}), b\left(\boldsymbol{k}^{\prime}\right)\right\}_{+}\right] \\
& =e \sum_{k}\left(N_{+}(\boldsymbol{k})-N_{-}(\boldsymbol{k})\right), \tag{9.56}
\end{align*}
$$

where

$$
N_{+}(\boldsymbol{k})=a^{\dagger}(\boldsymbol{k}) a(\boldsymbol{k}) \quad \text { and } \quad N_{-}(\boldsymbol{k})=b^{\dagger}(\boldsymbol{k}) b(\boldsymbol{k}) .
$$

Note that $Q=0$ if anticommutation relations (Fermi-Dirac statistics) are used.

The total charge of the field is

$$
\begin{equation*}
Q=e \sum_{k}\left(n_{+}(\boldsymbol{k})-n_{-}(\boldsymbol{k})\right), \tag{9.57}
\end{equation*}
$$

with

$$
n_{ \pm}(\boldsymbol{k})=0,1,2, \cdots, \infty
$$

Thus the operators $a(\boldsymbol{k}), a^{\dagger}(\boldsymbol{k}), N_{+}(\boldsymbol{k})$ denote the annihilation, creation and number operator for particle of electric charge $e$ and the operators $b(\boldsymbol{k}), b^{\dagger}(\boldsymbol{k}), N_{-}(\boldsymbol{k})$ denote the corresponding operators for a particle of charge $-e$. The total Hamiltonian is given by

$$
\begin{align*}
H & =\frac{1}{2} \sum_{k} \omega_{k} \sum_{r=1}^{2}\left\{a_{r}(\boldsymbol{k}), a_{r}^{\dagger}(\boldsymbol{k})\right\}_{+} \\
& \equiv \frac{1}{2} \sum_{k} \omega_{k}\left[\left\{a(\boldsymbol{k}), a^{\dagger}(\boldsymbol{k})\right\}_{+}+\left\{b(\boldsymbol{k}), b^{\dagger}(\boldsymbol{k})\right\}_{+}\right] \\
& =\sum_{k}\left[N_{+}(\boldsymbol{k})+N_{-}(\boldsymbol{k})\right] \omega_{k}+H_{0}, \tag{9.58}
\end{align*}
$$

with

$$
H_{0}=\sum_{k} \omega_{k} I
$$

The total field energy of the K-G field is

$$
\begin{equation*}
E_{K G}=E_{+}+E_{-}, \tag{9.59}
\end{equation*}
$$

where

$$
E_{ \pm}=\sum_{k}\left(n_{ \pm}(\boldsymbol{k})+\frac{1}{2}\right) \omega_{k}
$$

The total field energy would be zero if anticommutation relations are used for annihilation and creation operators. This means that the KleinGordon equation can be consistently quantized only by using the commutation relations corresponding to Bose-Einstein statistics.

### 9.2.2 Particles and antiparticles

From Eq. (9.57), we find that the complex scalar field indicates that there are two types of particles, one with charge $+e$ and the other with charge $-e$. Apart from the sign of the charge, both the particles have identical mass and other properties. So, they have to be associated with the particle and its antiparticle. This is a fundamental feature of relativistic quantum field theory, which is fully supported by experiments.

As an example of particle-antiparticle pair, we can cite $\pi^{+}$and $\pi^{-}$ mesons. They can be described by a complex Klein-Gordon field. For a real field, the charge operator $Q$ is identically zero and such a field can describe $\pi^{0}$ meson.

It is possible to extend the above consideration to other types of charges, for instance, hypercharge. The invariance of Lagrangian density $\mathscr{L}$ under phase transformation similar to (9.53) will allow the conservation of hypercharge. This will allow the occurrence of particles and antiparticles differing from each other in the sign of hypercharge. So, even electrically neutral particle can have a distinguishable antiparticle as in the case of neutral $K^{0}$ and $\bar{K}^{0}$. They have opposite hypercharge $Y= \pm 1$ and they can be represented by a complex Klein-Gordon field. Although electric charge is conserved in all types of interaction, hypercharge is conserved in strong and electromagnetic interactions but violated in weak interactions.

### 9.3 The covariant commutation relations

In relativistic field theories, calculations can be performed in such a way that the relativistic invariance of the theory can be exhibited at every stage. This is made possible by the introduction of the Lorentz invariant commutation relations.

The scalar field function $\Phi(\mathbf{x})$ can be written in a Lorentz invariant way as

$$
\begin{align*}
\Phi(\mathbf{x}) & =\frac{1}{(2 \pi)^{3 / 2}} \int_{k_{0}>0} \frac{d^{3} k}{\sqrt{2 \omega_{k}}}\left(a(\boldsymbol{k}) e^{-i \mathbf{k} \cdot \mathbf{x}}+a^{\dagger}(\boldsymbol{k}) e^{+i \mathbf{k} \cdot \mathbf{x}}\right) \\
& =\Phi^{+}(\mathbf{x})+\Phi^{-}(\mathbf{x}) \tag{9.60}
\end{align*}
$$

The destruction operator $a(\boldsymbol{k})$ has positive frequency dependence and the creation operator $a^{\dagger}(\boldsymbol{k})$ has negative frequency dependence. This distinction is Lorentz invariant and allows an invariant decomposition of $\Phi(\mathbf{x})$ into positive and negative frequency parts.

$$
\begin{align*}
\Phi^{+}(\mathbf{x}) & =\frac{1}{(2 \pi)^{3 / 2}} \int_{k_{0}>0} \frac{d^{3} k}{\sqrt{2 \omega_{k}}} a(\boldsymbol{k}) e^{-i \mathbf{k} \cdot \mathbf{x}}  \tag{9.61}\\
\Phi^{-}(\mathbf{x}) & =\frac{1}{(2 \pi)^{3 / 2}} \int_{k_{0}>0} \frac{d^{3} k}{\sqrt{2 \omega_{k}}} a^{\dagger}(\boldsymbol{k}) e^{i \mathbf{k} \cdot \mathbf{x}} \tag{9.62}
\end{align*}
$$

It can be verified that

$$
\left\{\Phi^{-}(\mathbf{x})\right\}^{\dagger}=\Phi^{+}(\mathbf{x})
$$

$\Phi^{+}(\mathbf{x})$ is a destruction operator and $\Phi^{-}(\mathbf{x})$ is a creation operator. The vacuum state $|0\rangle$ which has been characterized by $a(\boldsymbol{k})|0\rangle=0$ for all $\boldsymbol{k}$ has now the property

$$
\Phi^{+}(\mathbf{x})|0\rangle=0 .
$$

Let us now evaluate the commutator $\left[\Phi(\mathbf{x}), \Phi\left(\mathbf{x}^{\prime}\right)\right]_{-}$.

$$
\begin{align*}
& {\left[\Phi(\mathbf{x}), \Phi\left(\mathbf{x}^{\prime}\right)\right]_{-}} \\
& ==\Phi(\mathbf{x}) \Phi\left(\mathbf{x}^{\prime}\right)-\Phi\left(\mathbf{x}^{\prime}\right) \Phi(\mathbf{x}) \\
& =\frac{1}{(2 \pi)^{3 / 2}} \int_{k_{0}>0} \frac{d^{3} k}{\sqrt{2 \omega_{k}}} \int_{k_{0}^{\prime}>0} \frac{d^{3} k^{\prime}}{\sqrt{2 \omega_{k^{\prime}}}} \\
& \quad \times\left\{\left[a(\boldsymbol{k}), a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-} e^{i \mathbf{k}^{\prime} \cdot \mathbf{x}^{\prime}} e^{-i \mathbf{k} \cdot \mathbf{x}}+\left[a^{\dagger}(\boldsymbol{k}), a\left(\boldsymbol{k}^{\prime}\right)\right]_{-} e^{-i \mathbf{k}^{\prime} \cdot \mathbf{x}^{\prime}} e^{i \mathbf{k} \cdot \mathbf{x}}\right\} . \tag{9.63}
\end{align*}
$$

Using the commutation relations between the annihilation and creation operators

$$
\begin{aligned}
{\left[a(\boldsymbol{k}), a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-} } & =\delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right) \\
{\left[a(\boldsymbol{k}), a\left(\boldsymbol{k}^{\prime}\right)\right]_{-} } & =\left[a^{\dagger}(\boldsymbol{k}), a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=0,
\end{aligned}
$$

and integrating over $d^{3} k^{\prime}$, we get after replacing $\omega_{k}$ by $k_{0}$ (since $\omega_{k}^{2}=$ $\left.k_{0}^{2}=k^{2}+m^{2}\right)$

$$
\begin{align*}
{\left[\Phi(\mathbf{x}), \Phi\left(\mathbf{x}^{\prime}\right)\right]_{-} } & =\frac{1}{2(2 \pi)^{3}} \int_{k_{0}>0} \frac{d^{3} k}{k_{0}}\left(e^{-i \mathbf{k} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}-e^{i \mathbf{k} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}\right) \\
& =i \Delta\left(\mathbf{x}-\mathbf{x}^{\prime}\right), \tag{9.64}
\end{align*}
$$

with

$$
\begin{equation*}
\Delta\left(\mathbf{x}-\mathbf{x}^{\prime}\right)=-\frac{i}{2(2 \pi)^{3}} \int_{k_{0}>0} \frac{d^{3} k}{k_{0}}\left(e^{-i \mathbf{k} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}-e^{i \mathbf{k} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}\right) . \tag{9.65}
\end{equation*}
$$

Since

$$
\left[\Phi(\mathbf{x}), \Phi\left(\mathbf{x}^{\prime}\right)\right]_{-}=-\left[\Phi\left(\mathbf{x}^{\prime}\right), \Phi(\mathbf{x})\right]_{-}
$$

we get

$$
\begin{equation*}
\Delta\left(\mathrm{x}-\mathrm{x}^{\prime}\right)=-\Delta\left(\mathrm{x}^{\prime}-\mathrm{x}\right) \tag{9.66}
\end{equation*}
$$

It can be shown that $\Delta\left(x-x^{\prime}\right)$ satisfies the Klein-Gordon equation.

$$
\begin{equation*}
\left(\square_{\mathbf{x}}+\kappa^{2}\right) \Delta\left(\mathrm{x}-\mathrm{x}^{\prime}\right)=0 . \tag{9.67}
\end{equation*}
$$

The two-point function $\Delta\left(x-x^{\prime}\right)$ can be considered as the propagator of the scalar field.

It can be shown that

$$
\begin{equation*}
\left[\Phi^{+}(\mathbf{x}), \Phi^{+}\left(\mathbf{x}^{\prime}\right)\right]_{-}=\left[\Phi^{-}(\mathbf{x}), \Phi^{-}\left(\mathbf{x}^{\prime}\right)\right]_{-}=0 . \tag{9.68}
\end{equation*}
$$

In a similar way, we can also find the commutators $\left[\Phi^{+}(\mathbf{x}), \Phi^{-}\left(\mathbf{x}^{\prime}\right)\right]$ - and $\left[\Phi^{-}(\mathbf{x}), \Phi^{+}\left(\mathbf{x}^{\prime}\right)\right]_{-}$.

$$
\begin{align*}
& {\left[\Phi^{+}(\mathbf{x}), \Phi^{-}(\mathbf{x})\right]_{-}} \\
& =\Phi^{+}(\mathbf{x}) \Phi^{-}(\mathbf{x})-\Phi^{-}(\mathbf{x}) \Phi^{+}(\mathbf{x}) \\
& =\frac{1}{(2 \pi)^{3}} \int_{k_{0}>0} \frac{d^{3} k}{\sqrt{2 \omega_{k}}} \int_{k_{0}^{\prime}>0} \frac{d^{3} k^{\prime}}{\sqrt{2 \omega_{k^{\prime}}}} \\
& \times\left\{a(\boldsymbol{k}) a^{\dagger}\left(\boldsymbol{k}^{\prime}\right) e^{-i \mathbf{k} \cdot \mathbf{x}} e^{i \mathbf{k}^{\prime} \cdot \mathbf{x}^{\prime}}-a^{\dagger}\left(\boldsymbol{k}^{\prime}\right) a(\boldsymbol{k}) e^{i \mathbf{k}^{\prime} \cdot \mathbf{x}^{\prime}} e^{-i \mathbf{k} \cdot \mathbf{x}}\right\} . \tag{9.69}
\end{align*}
$$

Using the commutation relation

$$
\left[a(\boldsymbol{k}), a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=\delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right)
$$

and integrating over $d^{3} k^{\prime}$, we get after replacing $\omega_{k}$ by $k_{0}$

$$
\begin{align*}
{\left[\Phi^{+}(\mathbf{x}), \Phi^{-}(\mathbf{x})\right]_{-} } & =\frac{1}{(2 \pi)^{3}} \int_{k_{0}>0} \frac{d^{3} k}{2 k_{0}} e^{-i \mathbf{k} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \\
& =i \Delta^{+}\left(\mathbf{x}-\mathbf{x}^{\prime}\right) \tag{9.70}
\end{align*}
$$

In a similar way, we get

$$
\begin{align*}
{\left[\Phi^{-}(\mathbf{x}), \Phi^{+}\left(\mathbf{x}^{\prime}\right)\right]_{-} } & =-\frac{1}{(2 \pi)^{3}} \int_{k_{0}>0} \frac{d^{3} k}{2 k_{0}} e^{i \mathbf{k} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \\
& =i \Delta^{-}\left(\mathbf{x}-\mathbf{x}^{\prime}\right) \tag{9.71}
\end{align*}
$$

Thus, we find

$$
\begin{equation*}
\Delta^{+}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)+\Delta^{-}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)=\Delta\left(\mathrm{x}-\mathrm{x}^{\prime}\right) \tag{9.72}
\end{equation*}
$$

Likewise,

$$
\begin{align*}
& \Delta^{+}\left(\mathbf{x}-\mathbf{x}^{\prime}\right)-\Delta^{-}\left(\mathbf{x}-\mathbf{x}^{\prime}\right) \\
& \quad=-\frac{i}{(2 \pi)^{3}} \int_{k_{0}>0} \frac{d^{3} k}{2 k_{0}}\left\{e^{-i \mathbf{k} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}+e^{+i \mathbf{k} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}\right\} \\
& \quad=-\frac{i}{(2 \pi)^{3}} \int_{k_{0}>0} \frac{d^{3} k}{2 k_{0}} e^{i \boldsymbol{k} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)}\left\{e^{-i k_{0}\left(x_{0}-x_{0}^{\prime}\right)}+e^{+i k_{0}\left(x_{0}-x_{0}^{\prime}\right)}\right\} \\
& \quad=-\frac{i}{(2 \pi)^{3}} \int_{k_{0}>0} \frac{d^{3} k}{2 k_{0}} e^{i \boldsymbol{k} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)} 2 \cos k_{0}\left(x_{0}-x_{0}^{\prime}\right) \\
& \quad=-i \Delta^{(1)}\left(\mathbf{x}-\mathbf{x}^{\prime}\right) . \tag{9.73}
\end{align*}
$$

From Eqs. (9.72) and (9.73), it follows that

$$
\begin{align*}
& \Delta^{+}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)=\frac{1}{2}\left\{\Delta\left(\mathrm{x}-\mathrm{x}^{\prime}\right)-i \Delta^{(1)}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)\right\}  \tag{9.74}\\
& \Delta^{-}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)=\frac{1}{2}\left\{\Delta\left(\mathrm{x}-\mathrm{x}^{\prime}\right)+i \Delta^{(1)}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)\right\} \tag{9.75}
\end{align*}
$$

The covariant commutation relations play a crucial role in treating physical problems and performing calculations in a relativistically invariant way.

## Review Questions

9.1 Write dowm the Klein-Gordon equation for a real scalar field and discuss how it can be Fourier analyzed and quantized to show that it represents an assembly of Bose particles.
9.2 Show that the ground state of the system that the Klein-Gordon field represents is of infinite energy and explain how this difficulty can be overcome by normal ordering.
9.3 Show how the complex scalar field represents an assembly of charged particles whereas the real scalar field represents an assembly of uncharged particles.
9.4 Find the following covariant commutation relations between the scalar field functions at two different space-time points $\mathbf{x}$ and $\mathbf{x}^{\prime}$ :

$$
\left[\Phi(\mathbf{x}), \Phi\left(\mathbf{x}^{\prime}\right)\right]_{-}, \quad\left[\Phi^{+}(\mathbf{x}), \Phi^{-}\left(\mathbf{x}^{\prime}\right)\right]_{-}, \quad\left[\Phi^{-}(\mathbf{x}), \Phi^{+}\left(\mathbf{x}^{\prime}\right)\right]_{-}
$$

where $\Phi^{+}$and $\Phi^{-}$denote the positive and negative frequency parts of the field functions. Explain their physical significance.

## Problems

9.1 Given the Lagrangian density

$$
\mathscr{L}_{K G}(\mathbf{x})=\frac{1}{2}\left\{\sum_{\mu}\left(\frac{\partial \Phi}{\partial x_{\mu}}\right)^{2}-m^{2} \Phi^{2}\right\}
$$

deduce the Klein-Gordon field equation.
9.2 Given the field function $\Phi(\boldsymbol{x})$ and the corresponding conjugate momentum function $\Pi(\boldsymbol{x})$ for the Klein-Gordon field

$$
\begin{aligned}
\Phi(\boldsymbol{x}) & =\frac{1}{(2 \pi)^{3 / 2}} \int \frac{d^{3} k}{\sqrt{2 \omega_{k}}}\left(a(\boldsymbol{k}) e^{i \boldsymbol{k} \cdot \boldsymbol{x}}+a^{\dagger}(\boldsymbol{k}) e^{-i \boldsymbol{k} \cdot \boldsymbol{x}}\right) \\
\Pi(\boldsymbol{x}) & =\frac{1}{(2 \pi)^{3 / 2}} \int \frac{d^{3} k}{\sqrt{2 \omega_{k}}}\left(-i \omega_{k}\right)\left(a(\boldsymbol{k}) e^{i \boldsymbol{k} \cdot \boldsymbol{x}}-a^{\dagger}(\boldsymbol{k}) e^{-i \boldsymbol{k} \cdot \boldsymbol{x}}\right)
\end{aligned}
$$

obtain, by Fourier inversion, expressions for annihilation and creation operators $a(\boldsymbol{k})$ and $a^{\dagger}(\boldsymbol{k})$.
9.3 Show that the equal-time commutation relations

$$
\begin{aligned}
& {\left[\Phi(\mathbf{x}), \Pi\left(\mathbf{x}^{\prime}\right)\right]_{-}=i \hbar \delta\left(\mathbf{x}-\mathbf{x}^{\prime}\right)=i \delta\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \\
& {\left[\Phi(\mathbf{x}), \Phi\left(\mathbf{x}^{\prime}\right)\right]_{-}=\left[\Pi(\mathbf{x}), \Pi\left(\mathbf{x}^{\prime}\right)\right]_{-}=0}
\end{aligned}
$$

yield the commutation relations between the annihilation and creation operators.

$$
\begin{aligned}
{\left[a(\boldsymbol{k}), a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-} } & =\delta_{k, k^{\prime}} \\
{\left[a(\boldsymbol{k}), a\left(\boldsymbol{k}^{\prime}\right)\right]_{-} } & =\left[a^{\dagger}(\boldsymbol{k}), a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=0
\end{aligned}
$$

9.4 It is given that the prescription for going from discrete momentum to continuous momentum in the Fourier expansion of the field operator is

$$
\frac{1}{\sqrt{V}} \sum_{k} \longrightarrow \frac{1}{(2 \pi)^{3 / 2}} \int d^{3} k
$$

Deduce the above prescription.
9.5 Show that the covariant commutation relation given by Eq. (9.64) vanishes for equal times. Explain its physical significance.
9.6 Starting from the covariant commutation relation (9.64), deduce the canonical commutation relation for equal times:

$$
\left[\Pi(\boldsymbol{x}, t), \phi\left(\boldsymbol{x}^{\prime}, t\right)\right]_{-}=-i \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)
$$

9.7 Demonstrate that the field function $\Phi(\mathbf{x})$ given by Eq. (9.60) is really invariant under Lorentz transformation and the splitting of $\Phi(\mathbf{x})$ into positive and negative frequency parts $\Phi^{+}(\mathbf{x})$ and $\Phi^{-}(\mathbf{x})$ are also Lorentz invariant.
9.8 Given the Hamiltonian density $\mathscr{H}$ of the scalar field (Eq. (9.16),

$$
\mathscr{H}(\boldsymbol{x})=\frac{1}{2}\left(\Pi^{2}+(\boldsymbol{\nabla} \Phi)^{2}+m^{2} \Phi^{2}\right)
$$

deduce the expression (9.32) for the total Hamiltonian $H$.

$$
H=\int \mathscr{H}(\boldsymbol{x}) d^{3} x=\frac{1}{2} \sum_{k} \omega_{k}\left\{a^{\dagger}(\boldsymbol{k}) a(\boldsymbol{k})+a(\boldsymbol{k}) a^{\dagger}(\boldsymbol{k})\right\} .
$$

## Solutions to Problems

9.1 Given the Lagrangian density of the field, the field should obey the EulerLagrange equation

$$
\begin{equation*}
\frac{\partial \mathscr{L}}{\partial \Phi}-\frac{\partial}{\partial x^{\mu}} \frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \Phi}{\partial x^{\mu}}\right)}=0 \tag{9.76}
\end{equation*}
$$

The first term yields

$$
\begin{equation*}
\frac{\partial \mathscr{L}}{\partial \Phi}=-m^{2} \Phi \tag{9.77}
\end{equation*}
$$

The second term gives

$$
\begin{equation*}
\frac{\partial}{\partial x^{\mu}} \frac{\partial \mathscr{L}}{\partial\left(\frac{\partial \Phi}{\partial x^{\mu}}\right)}=\frac{\partial}{\partial x^{\mu}}\left\{\frac{\partial \Phi}{\partial x^{\mu}}\right\}=\frac{\partial^{2} \Phi}{\partial x^{\mu} \partial x^{\mu}} \tag{9.78}
\end{equation*}
$$

Substituting these in Euler-Lagrange equation, we get

$$
\begin{equation*}
\sum_{\mu} \frac{\partial^{2} \Phi}{\partial x^{\mu} \partial x^{\mu}}+m^{2} \Phi=0 \tag{9.79}
\end{equation*}
$$

This is the Klein-Gordon equation

$$
\begin{equation*}
\left(\square+m^{2}\right) \Phi=0 \tag{9.80}
\end{equation*}
$$

with

$$
\square=\frac{\partial^{2}}{\partial x^{0^{2}}}-\frac{\partial^{2}}{\partial x^{1^{2}}}-\frac{\partial^{2}}{\partial x^{2^{2}}}-\frac{\partial^{2}}{\partial x^{3^{2}}}=\frac{1}{c^{2}} \frac{\partial^{2}}{\partial t^{2}}-\nabla^{2}
$$

9.2 Multiply $\Phi(\boldsymbol{x})$ by $\omega_{k}$ and $\Pi(\boldsymbol{x})$ by $i$ and then adding and subtracting the two expressions, we get

$$
\begin{align*}
\omega_{k} \Phi(\boldsymbol{x})+i \Pi(\boldsymbol{x}) & =\frac{1}{(2 \pi)^{3 / 2}} \int \frac{d^{3} k}{\sqrt{2 \omega_{k}}} 2 \omega_{k} a(\boldsymbol{k}) e^{i \boldsymbol{k} \cdot \boldsymbol{x}}  \tag{9.81}\\
\omega_{k} \Phi(\boldsymbol{x})-i \Pi(\boldsymbol{x}) & =\frac{1}{(2 \pi)^{3 / 2}} \int \frac{d^{3} k}{\sqrt{2 \omega_{k}}} 2 \omega_{k} a^{\dagger}(\boldsymbol{k}) e^{-i \boldsymbol{k} \cdot \boldsymbol{x}} \tag{9.82}
\end{align*}
$$

Taking the Fourier transforms of (9.81) and (9.82), we get

$$
\begin{align*}
a(\boldsymbol{k}) & =\frac{1}{(2 \pi)^{3 / 2}} \int \frac{d^{3} x}{\sqrt{2 \omega_{k}}} e^{-i \boldsymbol{k} \cdot \boldsymbol{x}}\left\{\omega_{k} \Phi(\boldsymbol{x})+i \Pi(\boldsymbol{x})\right\} .  \tag{9.83}\\
a^{\dagger}(\boldsymbol{k}) & =\frac{1}{(2 \pi)^{3 / 2}} \int \frac{d^{3} x}{\sqrt{2 \omega_{k}}} e^{i \boldsymbol{k} \cdot \boldsymbol{x}}\left\{\omega_{k} \Phi(\boldsymbol{x})-i \Pi(\boldsymbol{x})\right\} \tag{9.84}
\end{align*}
$$

9.3 Using the expressions (9.83) and (9.84) for $a(\boldsymbol{k})$ and $a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)$, we get

$$
\begin{align*}
a(\boldsymbol{k}) a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)= & \frac{1}{(2 \pi)^{3}} \int \frac{d^{3} x}{\sqrt{2 \omega_{k}}} \frac{d^{3} x^{\prime}}{\sqrt{2 \omega_{k^{\prime}}}} e^{i\left(\boldsymbol{k}^{\prime} \cdot \boldsymbol{x}^{\prime}-\boldsymbol{k} \cdot \boldsymbol{x}\right)} \\
& \left\{\omega_{k} \Phi(\boldsymbol{x})+i \Pi(\boldsymbol{x})\right\}\left\{\omega_{k}^{\prime} \Phi\left(\boldsymbol{x}^{\prime}\right)-i \Pi\left(\boldsymbol{x}^{\prime}\right)\right\}  \tag{9.85}\\
a^{\dagger}\left(\boldsymbol{k}^{\prime}\right) a(\boldsymbol{k})= & \frac{1}{(2 \pi)^{3}} \int \frac{d^{3} x}{\sqrt{2 \omega_{k}}} \frac{d^{3} x^{\prime}}{\sqrt{2 \omega_{k^{\prime}}}} e^{i\left(\boldsymbol{k}^{\prime} \cdot \boldsymbol{x}^{\prime}-\boldsymbol{k} \cdot \boldsymbol{x}\right)} \\
& \left\{\omega_{k^{\prime}} \Phi\left(\boldsymbol{x}^{\prime}\right)-i \Pi\left(\boldsymbol{x}^{\prime}\right)\right\}\left\{\omega_{k} \Phi(\boldsymbol{x})+i \Pi(\boldsymbol{x})\right\} . \tag{9.86}
\end{align*}
$$

The commutator

$$
\begin{align*}
{\left[a(\boldsymbol{k}) a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=} & \frac{1}{(2 \pi)^{3}} \int \frac{d^{3} x}{\sqrt{2 \omega_{k}}} \frac{d^{3} x^{\prime}}{\sqrt{2 \omega_{k^{\prime}}}} e^{i\left(\boldsymbol{k}^{\prime} \cdot \boldsymbol{x}^{\prime}-\boldsymbol{k} \cdot \boldsymbol{x}\right)} \\
& \times\left\{\omega_{k} \omega_{k^{\prime}}\left[\Phi(\boldsymbol{x}), \Phi\left(\boldsymbol{x}^{\prime}\right)\right]_{-}+\left[\Pi(\boldsymbol{x}), \Pi\left(\boldsymbol{x}^{\prime}\right)\right]_{-}\right. \\
& \left.-i \omega_{k}\left[\Phi(\boldsymbol{x}), \Pi\left(\boldsymbol{x}^{\prime}\right)\right]_{-}+i \omega_{k^{\prime}}\left[\Pi(\boldsymbol{x}), \Phi\left(\boldsymbol{x}^{\prime}\right)\right]_{-}\right\}
\end{align*}
$$

Substituting the commutation relations between the field functions

$$
[\Phi(\boldsymbol{x}), \Pi(\boldsymbol{x})]_{-}=i \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right), \quad\left[\Phi(\boldsymbol{x}), \Phi\left(\boldsymbol{x}^{\prime}\right)\right]_{-}=\left[\Pi(\boldsymbol{x}), \Pi\left(\boldsymbol{x}^{\prime}\right)\right]_{-}=0
$$

The quantity within the curly bracket in Eq. (9.87) becomes

$$
\{\cdots\}=\omega_{k} \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)+\omega_{k^{\prime}} \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)
$$

The integration over $d^{3} x^{\prime}$ is trivial and is done by replacing $\boldsymbol{x}^{\prime}$ by $\boldsymbol{x}$ in the integrand.

$$
\begin{align*}
{\left[a(\boldsymbol{k}) a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-} } & =\frac{1}{(2 \pi)^{3}} \int \frac{d^{3} x}{\sqrt{2 \omega_{k}} \sqrt{2 \omega_{k^{\prime}}}} e^{i\left(\boldsymbol{k}^{\prime}-\boldsymbol{k}\right) \cdot \boldsymbol{x}}\left\{\omega_{k}+\omega_{k^{\prime}}\right\} \\
& =\delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right) \tag{9.88}
\end{align*}
$$

since

$$
\frac{1}{(2 \pi)^{3}} \int d^{3} x e^{i\left(\boldsymbol{k}^{\prime}-\boldsymbol{k}\right) \cdot \boldsymbol{x}}=\delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right)
$$

9.4 The number of states available with momentum lying between $k$ and $k+d k$ is

$$
\frac{V d^{3} k}{h^{3}}=\frac{V d^{3} k}{(2 \pi \hbar)^{3}}
$$

where $V$ is the normalization volume. The total number of states available in the case of momentum becoming a continuous variable is

$$
\int \frac{V d^{3} k}{(2 \pi \hbar)^{3}}
$$

So, the summation over the discrete number of states goes over into an integral if the momentum becomes a continuous variable.

$$
\begin{aligned}
\sum_{k} & \longrightarrow \int \frac{V d^{3} k}{(2 \pi \hbar)^{3}} \\
\frac{1}{V} \sum_{k} & \longrightarrow \int \frac{d^{3} k}{(2 \pi \hbar)^{3}} \text { in natural units } \int \frac{d^{3} k}{(2 \pi)^{3}}
\end{aligned}
$$

If the nrmalization volume is $V$, we usually take the normaliztion factor as $1 / \sqrt{V}$. Hence

$$
\frac{1}{\sqrt{V}} \sum_{k} \longrightarrow \frac{1}{(2 \pi)^{3 / 2}} \int d^{3} k
$$

9.5 The commutator can be expanded as

$$
\begin{aligned}
& {\left[\Phi(\mathbf{x}), \Phi\left(\mathbf{x}^{\prime}\right)\right]_{-}=\Phi(\mathbf{x}) \Phi\left(\mathbf{x}^{\prime}\right)-\Phi\left(\mathbf{x}^{\prime}\right) \Phi(\mathbf{x})} \\
& = \\
& \quad \frac{1}{(2 \pi)^{3 / 2}} \int_{k_{0}>0} \frac{d^{3} k}{\sqrt{2 \omega_{k}}} \int_{k_{0}^{\prime}>0} \frac{d^{3} k^{\prime}}{\sqrt{2 \omega_{k^{\prime}}}} \\
& \quad \times\left\{\left[a(\boldsymbol{k}), a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-} e^{i \mathbf{k}^{\prime} \cdot \mathbf{x}^{\prime}} e^{-i \mathbf{k} \cdot \mathbf{x}}+\left[a^{\dagger}(\boldsymbol{k}), a\left(\boldsymbol{k}^{\prime}\right)\right]_{-} e^{-i \mathbf{k}^{\prime} \cdot \mathbf{x}^{\prime}} e^{i \mathbf{k} \cdot \mathbf{x}}\right\}
\end{aligned}
$$

Using the commutation relations between the annihilation and creation operators

$$
\begin{aligned}
{\left[a(\boldsymbol{k}), a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-} } & =\delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right) \\
{\left[a(\boldsymbol{k}), a\left(\boldsymbol{k}^{\prime}\right)\right]_{-} } & =\left[a^{\dagger}(\boldsymbol{k}), a^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=0
\end{aligned}
$$

and integrating over $d^{3} k^{\prime}$, we get

$$
\begin{aligned}
{\left[\Phi(\mathbf{x}), \Phi\left(\mathbf{x}^{\prime}\right)\right]_{-} } & =\frac{1}{2(2 \pi)^{3}} \int_{k_{0}>0} \frac{d^{3} k}{k_{0}}\left(e^{-i \mathbf{k} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}-e^{i \mathbf{k} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}\right) \\
& =i \Delta\left(\mathbf{x}-\mathbf{x}^{\prime}\right)
\end{aligned}
$$

with

$$
\begin{aligned}
\Delta\left(\mathbf{x}-\mathbf{x}^{\prime}\right) & =-\frac{i}{2(2 \pi)^{3}} \int_{k_{0}>0} \frac{d^{3} k}{k_{0}}\left(e^{-i \mathbf{k} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}-e^{i \mathbf{k} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}\right) \\
& =\frac{1}{(2 \pi)^{3}} \int_{k_{0}>0} \frac{d^{3} k}{k_{0}} e^{i \boldsymbol{k} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)} \sin k_{0}\left(x_{0}-x_{0}^{\prime}\right) \\
& =0, \quad \text { if } x_{0}=x_{0}^{\prime}
\end{aligned}
$$

Hence the commutator $\left[\Phi(\mathbf{x}), \Phi\left(\mathbf{x}^{\prime}\right)\right]_{-}$vanishes for equal times.
A vanishing commutator allows one to simultaneously diagonalize the two operators and permits to make precise observation simultaneously of the two observables corresponding to the operators. So, the measurement of the Hermitian field operators $\Phi(\boldsymbol{x}, t)$ and $\Phi\left(\boldsymbol{x}^{\prime}, t\right)$ cannot interfere with one another, since no light signal can connect these two space-time points. It follows that the commutator vanishes for any space-like intervals $i . e$ for space-time intervals $\boldsymbol{x}_{1}-\boldsymbol{x}_{2}>c\left(t_{1}-t_{2}\right)$.
9.6 Let us start with the commutation relation of the scalar field operators.

$$
\begin{aligned}
{\left[\Phi(\mathbf{x}), \Phi\left(\mathbf{x}^{\prime}\right)\right]_{-} } & =i \hbar \Delta\left(\mathbf{x}-\mathbf{x}^{\prime}\right) \\
\frac{\partial}{\partial t}\left[\Phi(\boldsymbol{x}, t), \Phi\left(\boldsymbol{x}^{\prime}, t^{\prime}\right)\right]_{-} & =\left[\Pi(\boldsymbol{x}, t), \Phi\left(\boldsymbol{x}^{\prime}, t^{\prime}\right)\right]_{-}
\end{aligned}
$$

since $\frac{\partial}{\partial t} \Phi\left(\boldsymbol{x}^{\prime}, t^{\prime}\right)=0$ and $\frac{\partial}{\partial t} \Phi(x, t)=\Pi$.
Hence

$$
\left[\Pi(\boldsymbol{x}, t), \Phi\left(\boldsymbol{x}^{\prime}, t^{\prime}\right)\right]_{-}=\frac{\partial}{\partial t}\left[\Phi(\boldsymbol{x}, t), \Phi\left(\boldsymbol{x}^{\prime}, t^{\prime}\right)\right]_{-}=i \hbar \frac{\partial}{\partial t} \Delta\left(\mathbf{x}-\mathbf{x}^{\prime}\right)
$$

where

$$
\begin{aligned}
\Delta\left(\mathbf{x}-\mathbf{x}^{\prime}\right) & =-\frac{i}{2} \frac{1}{(2 \pi)^{3}} \int_{k_{0}>0} \frac{d^{3} k}{k_{0}}\left(e^{-i \mathbf{k} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}-\left(e^{i \mathbf{k} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}\right)\right. \\
& =-\frac{i}{2} \frac{1}{(2 \pi)^{3}} \int_{k_{0}>0} \frac{d^{3} k}{k_{0}} e^{i \boldsymbol{k} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)}\left\{2 i \sin \left(k_{0}\left(x_{0}-x_{0}^{\prime}\right)\right)\right\}
\end{aligned}
$$

The only factors that depend on time are $x_{0}=c t$ and $x_{0}^{\prime}=c t^{\prime}$. So, differentiating with respect to $t$, we get

$$
\frac{\partial}{\partial t} \sin \left(k_{0}\left(x_{0}-x_{0}^{\prime}\right)\right)=k_{0} \cos \left(k_{0}\left(x_{0}-x_{0}^{\prime}\right)\right)=k_{0}, \quad \text { if } x_{0}=x_{0}^{\prime}
$$

Substituting this, we get

$$
\begin{aligned}
\Delta\left(\mathbf{x}-\mathbf{x}^{\prime}\right) & =\frac{1}{(2 \pi)^{3}} \int d^{3} k e^{i \boldsymbol{k} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)} \\
& =\delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)
\end{aligned}
$$

Thus, we get back the canonical commutation rule for equal times.

$$
\left[\Pi(\boldsymbol{x}, t), \Phi\left(\boldsymbol{x}^{\prime}, t\right)\right]_{-}=i \hbar \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)
$$

9.7 Expressions (9.60) - (9.62) for $\Phi(\mathbf{x}), \Phi^{+}(\mathbf{x}), \Phi^{-}(\mathbf{x})$, are not explicitly in Lorentz invariant forms since they involve three dimensional vectors but they can be cast in terms of four vectors using the following identity.

$$
\begin{aligned}
\Phi(\mathbf{x}) & =\frac{1}{(2 \pi)^{3 / 2}} \int_{k_{0}>0} \frac{d^{3} k}{\sqrt{2 \omega_{k}}}\left(a(\boldsymbol{k}) e^{-i \mathbf{k} \cdot \mathbf{x}}+a^{\dagger}(\boldsymbol{k}) e^{+i \mathbf{k} \cdot \mathbf{x}}\right) \\
& =\Phi^{+}(\mathbf{x})+\Phi^{-}(\mathbf{x})
\end{aligned}
$$

First let us show that

$$
\begin{aligned}
\frac{d^{4} k}{(2 \pi)^{4}} 2 \pi \delta\left(\mathbf{k}^{2}-m^{2}\right) \theta\left(k_{0}\right)= & \frac{d^{4} k}{(2 \pi)^{3}} \delta\left(k_{0}^{2}-\omega_{k}^{2}\right) \theta\left(k_{0}\right), \\
& \text { since } \mathbf{k}^{2}-m^{2}=k_{0}^{2}-\boldsymbol{k}^{2}-m^{2}=k_{0}^{2}-\omega^{2} \\
= & \frac{d^{4} k}{(2 \pi)^{3}} \delta\left[\left(k_{0}+\omega_{k}\right)\left(k_{0}-\omega_{k}\right)\right] \theta\left(k_{0}\right) \\
= & \frac{d^{4} k}{(2 \pi)^{3}} \frac{1}{2 k_{0}}\left[\delta\left(k_{0}+\omega_{k}\right)+\delta\left(k_{0}-\omega_{k}\right)\right] \theta\left(k_{0}\right), \\
& \text { since } \delta\left(x^{2}-a^{2}\right)=\frac{1}{2 a}\{\delta(x+a)+\delta(x-a)\} \\
= & \frac{d^{3} k}{(2 \pi)^{3}} \frac{d k_{0}}{2 k_{0}} \delta\left(k_{0}-\omega_{k}\right) \\
= & \frac{d^{3} k}{(2 \pi)^{3}} \frac{1}{2 \omega_{k}} .
\end{aligned}
$$

Using the above identity and adjusting the normalization factor, we can display the above quantities in a Lorentz invariant form.
9.8 The Hamiltonian density $\mathscr{H}$ (Eq. 9.16) consists of three terms. To find the total Hamiltonian, one has to perform integration over $d^{3} x$ on each
term. In the absence of interaction, the Hamiltonian will be constant in time. So, we can choose $t=0$.
Let us do the integration over the first term in some detail. Using Eq. (9.22), we get

$$
\begin{aligned}
& \int \Pi^{2}(\boldsymbol{x}) d^{3} x=\frac{1}{(2 \pi)^{3}} \int \frac{d^{3} k d^{3} k^{\prime} d^{3} x}{\sqrt{2 \omega_{k} 2 \omega_{k^{\prime}}}}\left(-\omega_{k} \omega_{k^{\prime}}\right) \\
& \quad \times\left\{a(\boldsymbol{k}) e^{i \boldsymbol{k} \cdot \boldsymbol{x}}-a^{\dagger}(\boldsymbol{k}) e^{-i \boldsymbol{k} \cdot \boldsymbol{x}}\right\}\left\{a\left(\boldsymbol{k}^{\prime}\right) e^{i \boldsymbol{k}^{\prime} \cdot \boldsymbol{x}}-a^{\dagger}\left(\boldsymbol{k}^{\prime}\right) e^{-i \boldsymbol{k}^{\prime} \cdot \boldsymbol{x}}\right\} .
\end{aligned}
$$

First let us perform integration over $d^{3} x$. The $\boldsymbol{x}$ - dependent terms are only within the curly brackets.

$$
\begin{aligned}
\int d^{3} x\{\cdots\}\{\cdots\}= & \int d^{3} x\left[a(\boldsymbol{k}) a\left(\boldsymbol{k}^{\prime}\right) e^{i\left(\boldsymbol{k}+\boldsymbol{k}^{\prime}\right) \cdot \boldsymbol{x}}-a(\boldsymbol{k}) a^{\dagger}\left(\boldsymbol{k}^{\prime}\right) e^{i\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right) \cdot \boldsymbol{x}}\right. \\
& \left.-a^{\dagger}(\boldsymbol{k}) a\left(\boldsymbol{k}^{\prime}\right) e^{-i\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right) \cdot \boldsymbol{x}}+a^{\dagger}(\boldsymbol{k}) a^{\dagger}\left(\boldsymbol{k}^{\prime}\right) e^{-i\left(\boldsymbol{k}+\boldsymbol{k}^{\prime}\right) \cdot \boldsymbol{x}}\right] \\
= & (2 \pi)^{3}\left[a(\boldsymbol{k}) a\left(\boldsymbol{k}^{\prime}\right) \delta\left(\boldsymbol{k}+\boldsymbol{k}^{\prime}\right)-a(\boldsymbol{k}) a^{\dagger}\left(\boldsymbol{k}^{\prime}\right) \delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right)\right. \\
& \left.-a^{\dagger}(\boldsymbol{k}) a\left(\boldsymbol{k}^{\prime}\right) \delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right)+a^{\dagger}(\boldsymbol{k}) a^{\dagger}\left(\boldsymbol{k}^{\prime}\right) \delta\left(\boldsymbol{k}+\boldsymbol{k}^{\prime}\right)\right]
\end{aligned}
$$

where we have used the standard integral

$$
\int e^{i\left(\boldsymbol{k}+\boldsymbol{k}^{\prime}\right) \cdot \boldsymbol{x}} d^{3} x=(2 \pi)^{3} \delta\left(\boldsymbol{k}+\boldsymbol{k}^{\prime}\right)
$$

The integration over $d^{3} k^{\prime}$ is trivial because of the delta function. Thus, we get

$$
\begin{align*}
\int \Pi^{2}(\boldsymbol{x}) d^{3} x= & \int d^{3} k\left(\frac{\omega_{k}}{2}\right)\left\{a(\boldsymbol{k}) a^{\dagger}(\boldsymbol{k})+a^{\dagger}(\boldsymbol{k}) a(\boldsymbol{k})\right. \\
& \left.-a(\boldsymbol{k}) a(-\boldsymbol{k})-a^{\dagger}(\boldsymbol{k}) a^{\dagger}(-\boldsymbol{k})\right\} \tag{9.89}
\end{align*}
$$

In the same way, the integrals over the other two terms can be evaluated and we give below only the final results.

$$
\begin{align*}
\int(\boldsymbol{\nabla} \Phi(\boldsymbol{x}))^{2} d^{3} x= & \int d^{3} k\left(\frac{k^{2}}{2 \omega_{k}}\right)\left\{a(\boldsymbol{k}) a^{\dagger}(\boldsymbol{k})+a^{\dagger}(\boldsymbol{k}) a(\boldsymbol{k})\right. \\
& \left.+a(\boldsymbol{k}) a(-\boldsymbol{k})+a^{\dagger}(\boldsymbol{k}) a^{\dagger}(-\boldsymbol{k})\right\} .  \tag{9.90}\\
m^{2} \int(\Phi(\boldsymbol{x}))^{2} d^{3} x= & \int d^{3} k\left(\frac{m^{2}}{2 \omega_{k}}\right)\left\{a(\boldsymbol{k}) a^{\dagger}(\boldsymbol{k})+a^{\dagger}(\boldsymbol{k}) a(\boldsymbol{k})\right. \\
& \left.+a(\boldsymbol{k}) a(-\boldsymbol{k})+a^{\dagger}(\boldsymbol{k}) a^{\dagger}(-\boldsymbol{k})\right\} . \tag{9.91}
\end{align*}
$$

Adding (9.89), (9.90) and (9.91) and observing that $\omega_{k}^{2}=k^{2}+m^{2}$, we get

$$
H=\frac{1}{2} \int \omega_{k}\left[a^{\dagger}(\boldsymbol{k}) a(\boldsymbol{k})+a(\boldsymbol{k}) a^{\dagger}(\boldsymbol{k})\right] .
$$

## Chapter 10

## The Dirac Field

In the last chapter, we have considered the scalar fields which, when quantized, using commutation relations yield an assembly of particles which obey Bose-Einstein statistics. In this chapter, let us consider the Dirac spin- $\frac{1}{2}$ field which is described by a four-component spinor $\psi$. Each component $\psi_{\alpha}$ of the four-component spinor should be treated as independent dynamical variable. The Dirac field, when quantized, yields an assembly of spin- $\frac{1}{2}$ particles known as Fermions. They obey the Fermi-Dirac statistics and satisfy the Pauli exclusion principle. This is achieved by replacing the commutation relations between the annihilation and creation operators by anti-commutation relations.

It is found that Bosons which have integral spins can be quantized only by invoking commutation relations and the Fermions which have half-integral spins can be quantized only by invoking anticommutation relations. The close relation between spin and statistics and their connection to commutation or anticommutation relations between annihilation and creation operators are the essential features of relativistic quantum field theory.

The reader is advised to read in advance Chap. 2, where the Dirac equation is treated as a single particle relativistic wave equation. This is a prerequisite for studying the Dirac equation as a field equation with a field function that can be expanded in terms of a complete set of eigenfunctions of the free particle Dirac equation. Anyhow, a brief summary of the relevant results are given for the benefit of the reader.

### 10.1 Plane wave solutions of the Dirac Equation

Let us consider the Dirac equation in Feynman's notation ${ }^{1}$ for a particle of mass $m$ (in natural units $\hbar=c=1$ ).

$$
\begin{equation*}
(\not p-m) \psi(\boldsymbol{x}, t)=0 \quad \text { or } \quad\left(i \gamma^{\mu} \partial_{\mu}-m\right) \psi(\mathbf{x})=0 . \tag{10.1}
\end{equation*}
$$

The symbol $\partial_{\mu}$ stands for the differential $\partial / \partial x^{\mu}$ and $\mathbf{x}$ denotes the fourvector $x_{0}=c t, x_{1}, x_{2}, x_{3}$. The repeated index $\mu$ implies summation over the index $\mu$. They are, in total, four solutions for Eq. (10.1). Two of them correspond to positive energy $E_{p}$, one with spin up and the other with spin down. The other two correspond to negative energy $-E_{p}$, one with spin up and the other with spin down. The four plane wave solutions of Eq. (10.1) are ${ }^{2}$

$$
\begin{align*}
\varphi_{r}(\mathbf{x}) & =\frac{1}{\sqrt{V}} \frac{1}{\sqrt{2 E_{p}}} u_{r}(\boldsymbol{p}) e^{-i \mathbf{p} \cdot \mathbf{x}}, & r=1,2  \tag{10.2}\\
\varphi_{r}(\mathbf{x}) & =\frac{1}{\sqrt{V}} \frac{1}{\sqrt{2 E_{p}}} u_{r}(\boldsymbol{p}) e^{i \mathbf{p} \cdot \mathbf{x}}, & r=3,4 \tag{10.3}
\end{align*}
$$

The scalar product of four vectors $\mathbf{p}$ and $\mathbf{x}$ are defined by

$$
\begin{equation*}
\mathbf{p} \cdot \mathbf{x}=E t-\boldsymbol{p} \cdot \boldsymbol{x} ; \quad \mathbf{p}^{2}=p^{\mu} p_{\mu}=E^{2}-\boldsymbol{p}^{2}=m^{2} \tag{10.4}
\end{equation*}
$$

We use the convention of denoting the four-vectors by upright bold letters $(\mathbf{p}, \mathbf{x})$, three-vectors by italic bold letters $(\boldsymbol{p}, \boldsymbol{x})$ and scalars by ordinary italics $(p, x)$. In Eqs. (10.2) and (10.3), $u_{r}(\boldsymbol{p})$ denotes the spinors which are really four-component but expressed below in the two-component form for the sake of compactness.

$$
\begin{array}{ll}
u_{r}(\boldsymbol{p})=\left(E_{p}+m\right)^{1 / 2}\left[\begin{array}{c}
\chi_{ \pm} \\
\frac{\boldsymbol{\sigma} \cdot \boldsymbol{p}}{E_{p}+m} \chi_{ \pm}
\end{array}\right], & r=1,2 \\
u_{r}(\boldsymbol{p})=\left(E_{p}+m\right)^{1 / 2}\left[\begin{array}{c}
\boldsymbol{\sigma} \cdot \boldsymbol{p} \\
E_{p}+m \\
\chi_{ \pm} \\
\chi_{ \pm}
\end{array}\right], & r=3,4 . \tag{10.6}
\end{array}
$$

$$
\begin{aligned}
& { }^{1} \text { The Feynman notation } p \text { stands for } \\
& \qquad \begin{aligned}
p & =\gamma^{\mu} p_{\mu}=\gamma^{0} E-\gamma \cdot \boldsymbol{p}=i\left(\gamma^{0} \frac{\partial}{\partial x^{0}}+\gamma^{1} \frac{\partial}{\partial x^{1}}+\gamma^{2} \frac{\partial}{\partial x^{2}}+\gamma^{3} \frac{\partial}{\partial x^{3}}\right) \\
& =i \sum_{\mu} \gamma^{\mu} \frac{\partial}{\partial x^{\mu}}=i \gamma^{\mu} \partial_{\mu} .
\end{aligned}
\end{aligned}
$$

${ }^{2}$ The suffix $r$ denotes the different spin orientations.

In Eqs. (10.5) and (10.6), the symbols $E_{p}, \chi_{+}$and $\chi_{-}$stand for

$$
E_{p}=+\sqrt{p^{2}+m^{2}} ; \quad \chi_{+}=\left[\begin{array}{l}
1  \tag{10.7}\\
0
\end{array}\right] ; \quad \chi_{-}=\left[\begin{array}{l}
0 \\
1
\end{array}\right]
$$

Different authors use different normalizations for the spinors. We choose the following normalization for the spinors $u_{r}(\boldsymbol{p})$, their Hermitian conjugates $u_{r}^{\dagger}(\boldsymbol{p})$ and their Hermitian adjoints $\bar{u}_{r}(\boldsymbol{p})=u_{r}^{\dagger}(\boldsymbol{p}) \gamma_{0}$, the normalization volume being $V$.

$$
\begin{array}{llrl}
u_{r}^{\dagger}(\boldsymbol{p}) u_{s}(\boldsymbol{p})=2 E_{p} \delta_{r s} ; & & r, s=1,2,3,4 . \\
\bar{u}_{r}(\boldsymbol{p}) u_{s}(\boldsymbol{p})=2 m \delta_{r s} ; & & r, s=1,2 . \\
\bar{u}_{r}(\boldsymbol{p}) u_{s}(\boldsymbol{p})=-2 m \delta_{r s} ; & & r, s=3,4 . \\
\bar{u}_{r}(\boldsymbol{p}) u_{s}(\boldsymbol{p})=0 & & r=1,2 ; s=3,4 . \\
\bar{u}_{r}(\boldsymbol{p}) u_{s}(\boldsymbol{p})=0 & & r=3,4 ; s=1,2 . \tag{10.12}
\end{array}
$$

The spinors satisfy the following equations:

$$
\begin{align*}
\left(\gamma^{\mu} p_{\mu}-m\right) u_{r}(\boldsymbol{p}) & =0, & & r=1,2 ;  \tag{10.13}\\
\left(\gamma^{\mu} p_{\mu}+m\right) u_{r}(\boldsymbol{p}) & =0, & & r=3,4 ; \tag{10.14}
\end{align*}
$$

which are obtained from Eq. (10.1) using the relation

$$
p_{\mu}=-i \hbar \frac{\partial}{\partial x^{\mu}} \xrightarrow[\text { units }]{\text { In natural }}-i \frac{\partial}{\partial x^{\mu}}=-i \partial_{\mu} .
$$

The spinor normalizations are such that

$$
\int_{V} \varphi_{r}^{\dagger}(\mathbf{x}) \varphi_{s}(\mathbf{x}) d^{3} x=2 E_{p} \delta_{r s}, \quad r, s=1,2,3,4
$$

The solutions $\varphi_{1}(\mathbf{x})$ and $\varphi_{2}(\mathbf{x})$ correspond to positive energy $\left(E=E_{p}\right)$ with momentum $\boldsymbol{p}$ whereas the solutions $\varphi_{3}(\mathbf{x})$ and $\varphi_{4}(\mathbf{x})$ correspond to negative energy ( $E=-E_{p}$ ) with momentum $-\boldsymbol{p}$.

The Hermitian adjoint of the Dirac equation (10.1) is

$$
\begin{equation*}
\bar{\psi}(\mathbf{x})\left(i \gamma^{\mu} \overleftarrow{\partial_{\mu}}+m\right)=0 \tag{10.15}
\end{equation*}
$$

with $\bar{\psi}(\mathbf{x})=\psi^{\dagger}(\mathbf{x}) \gamma_{0}$. The left-arrow sign indicates that the operator acts on the function located on the left.

### 10.2 Lagrangian density for the Dirac field

It can be easily verified that both the Dirac equation (10.1) and its Hermitian adjoint (10.15) can be obtained from the Lagrangian density

$$
\begin{equation*}
\mathscr{L}_{D}(\mathbf{x})=\bar{\psi}(\mathbf{x})\left(i \gamma^{\mu} \partial_{\mu}-m\right) \psi(\mathbf{x}) \tag{10.16}
\end{equation*}
$$

by insisting that it satisfies the Euler-Lagrange equation.
If $\psi$, is a complex field, then we can treat both $\psi$ and $\bar{\psi}$ as independent fields. The Euler-Lagrange equation for $\bar{\psi}$ gives

$$
\begin{equation*}
\partial_{\mu}\left(\frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \bar{\psi}\right)}\right)=\frac{\partial \mathscr{L}}{\partial \bar{\psi}} . \tag{10.17}
\end{equation*}
$$

Since the Lagrangian density $\mathscr{L}$ given by Eq. (10.16) does not contain any derivative of $\bar{\psi}$, the left hand side of Eq. (10.17) vanishes and Eq. (10.17) yields the Dirac equation (10.1).

$$
\frac{\partial \mathscr{L}}{\partial \bar{\psi}}=\left(i \gamma^{\mu} \partial_{\mu}-m\right) \psi(\mathbf{x})=0
$$

In a similar way, we can consider the Euler-Lagrange equation for $\psi$ and obtain the Hermitian conjugate of the Dirac equation (10.15).

In quantum field theory, the Lagrangian density should be Hermitian. In the case of scalar fields, we found the corresponding Lagrangian to be Hermitian. But the Dirac Lagrangian (10.16) is not Hermitian. The mass term is Hermitian whereas the Hermitian conjugate of the other term is given by

$$
\begin{equation*}
\left(i \bar{\psi} \gamma^{\mu} \partial_{\mu} \psi\right)^{\dagger}=-i\left(\partial_{\mu} \psi\right)^{\dagger} \gamma^{\mu \dagger} \gamma_{0} \psi=-i\left(\partial_{\mu} \bar{\psi}\right) \gamma^{\mu} \psi \tag{10.18}
\end{equation*}
$$

using the hermiticity property of the gamma matrices

$$
\begin{equation*}
\left(\gamma^{\mu}\right)^{\dagger}=\gamma^{0} \gamma^{\mu} \gamma^{0} \tag{10.19}
\end{equation*}
$$

Instead of (10.16), we can use the following Hermitian Lagrangian density:

$$
\begin{equation*}
\mathscr{L}^{\prime}=\frac{i}{2}\left(\bar{\psi} \gamma^{\mu} \partial_{\mu} \psi-\left(\partial_{\mu} \bar{\psi}\right) \gamma^{\mu} \psi\right)-m \bar{\psi} \psi \tag{10.20}
\end{equation*}
$$

However, it is not necessary since the two Lagrangians $\mathscr{L}$ and $\mathscr{L}^{\prime}$ are equivalent since they differ by a term which can be expressed as a total divergence. The difference

$$
\begin{equation*}
\mathscr{L}-\mathscr{L}^{\prime}=\frac{i}{2}\left(\bar{\psi} \gamma^{\mu} \partial_{\mu} \psi+\left(\partial_{\mu} \bar{\psi}\right) \gamma^{\mu} \psi\right)=\partial_{\mu}\left(\frac{i}{2} \bar{\psi} \gamma^{\mu} \psi\right) \tag{10.21}
\end{equation*}
$$

is a total divergence. It may be recalled that the two Lagrangians which differ by a total divergence would yield the same equations of motion based on the action principle. So, we can continue to use the same Lagrangian (10.16).

Now, let us proceed to find the conjugate momentum density and the Hamiltonian density of the field.

$$
\begin{align*}
\Pi(\mathbf{x}) & =\frac{\partial \mathscr{L}_{D}}{\partial \dot{\psi}}=\frac{\partial \mathscr{L}_{D}}{\partial \psi_{, 0}}=i \bar{\psi}(\mathbf{x}) \gamma_{0}=i \psi^{\dagger}(\mathbf{x})  \tag{10.22}\\
\mathscr{H}(\mathbf{x}) & =\Pi(\mathbf{x}) \dot{\psi}(\mathbf{x})-\mathscr{L}_{D}(\mathbf{x}) \\
& =\Pi(\mathbf{x}) \dot{\psi}(\mathbf{x})=i \psi^{\dagger}(\mathbf{x}) \dot{\psi}(\mathbf{x})=i \psi^{\dagger}(\mathbf{x}) \frac{\partial \psi(\mathbf{x})}{\partial t} \tag{10.23}
\end{align*}
$$

The last line in Eq. (10.23) is obtained using the result that the Lagrangian density $\mathscr{L}_{D}$ given by Eq. (10.16) vanishes if the field function satisfies the Dirac equation.

Our aim is to find the total Hamiltonian $H$.

$$
\begin{equation*}
H=\int \mathscr{H} d^{3} x=\int \psi^{\dagger}(\mathbf{x}) i \frac{\partial \psi(\mathbf{x})}{\partial t} d^{3} x \tag{10.24}
\end{equation*}
$$

### 10.3 Fourier decomposition of the Dirac field

Let us now replace the negative energy electron spinors $u_{3}(-\boldsymbol{p})$ and $u_{4}(-\boldsymbol{p})$ by positive energy positron spinors $v_{1}(\boldsymbol{p})$ and $v_{2}(\boldsymbol{p})$,

$$
\begin{equation*}
v_{1}(\boldsymbol{p}) \equiv u_{3}(-\boldsymbol{p}), \quad v_{2}(\boldsymbol{p}) \equiv u_{4}(-\boldsymbol{p}), \tag{10.25}
\end{equation*}
$$

and treat the Dirac field functions $\psi(\mathbf{x})$ and $\bar{\psi}(\mathbf{x})$ as field operators. The field operators $\psi(\mathbf{x})$ and $\bar{\psi}(\mathbf{x})$ are expanded in terms of $\varphi_{r}(\mathbf{x})$, defined by Eqs. (10.2) and (10.3).
$\psi(\mathbf{x})=\frac{1}{\sqrt{V}} \sum_{p} \sqrt{\frac{1}{2 E_{p}}} \sum_{r=1}^{2}\left[c_{r}(\boldsymbol{p}) u_{r}(\boldsymbol{p}) e^{-i \mathbf{p} \cdot \mathbf{x}}+d_{r}^{\dagger}(\boldsymbol{p}) v_{r}(\boldsymbol{p}) e^{i \mathbf{p} \cdot \mathbf{x}}\right](; 10.26)$
$\bar{\psi}(\mathbf{x})=\frac{1}{\sqrt{V}} \sum_{p} \sqrt{\frac{1}{2 E_{p}}} \sum_{r=1}^{2}\left[c_{r}^{\dagger}(\boldsymbol{p}) \bar{u}_{r}(\boldsymbol{p}) e^{i \mathbf{p} \cdot \mathbf{x}}+d_{r}(\boldsymbol{p}) \bar{v}_{r}(\boldsymbol{p}) e^{-i \mathbf{p} \cdot \mathbf{x}}\right]$.
The expansion coefficients $c_{r}(\boldsymbol{p})$ and $c_{r}^{\dagger}(\boldsymbol{p})$ are to be interpreted as annihilation and creation operators of the Fermions and the operators $d_{r}(\boldsymbol{p})$ and $d_{r}^{\dagger}(\boldsymbol{p})$ are to be interpreted as annihilation and creation operators
of the anti-Fermions. For this, let us write down $\psi^{\dagger}$ and $i \frac{\partial \psi}{\partial t}$, using the expansion (10.26) for $\psi(\mathbf{x})$.

$$
\begin{align*}
& \psi^{\dagger}(\mathbf{x})=\frac{1}{\sqrt{V}} \sum_{p} \sqrt{\frac{1}{2 E_{p}}} \sum_{r=1}^{2}\left[c_{r}^{\dagger}(\boldsymbol{p}) u_{r}^{\dagger}(\boldsymbol{p}) e^{i \mathbf{p} \cdot \mathbf{x}}+d_{r}(\boldsymbol{p}) v_{r}^{\dagger}(\boldsymbol{p}) e^{-i \mathbf{p} \cdot \mathbf{x}}\right]  \tag{10.28}\\
& i \frac{\partial \psi(\mathbf{x})}{\partial t}=\frac{1}{\sqrt{V}} \sum_{p^{\prime}} \sqrt{\frac{E_{p^{\prime}}}{2}} \sum_{s=1}^{2}\left[c_{s}\left(\boldsymbol{p}^{\prime}\right) u_{s}\left(\boldsymbol{p}^{\prime}\right) e^{-i \mathbf{p}^{\prime} \cdot \mathbf{x}}\right. \\
& \left.\left.-d_{s}^{\dagger}\left(\boldsymbol{p}^{\prime}\right) v_{s}\left(\boldsymbol{p}^{\prime}\right)\right) e^{i \mathbf{p}^{\prime} \cdot \mathbf{x}}\right] . \tag{10.29}
\end{align*}
$$

Substituting expressions (10.28) and (10.29) in Eq. (10.24), we get

$$
\begin{align*}
& H= \frac{1}{V} \sum_{p p^{\prime}} \\
& \sqrt{\frac{E_{p^{\prime}}}{4 E_{p}}} \int \sum_{r, s}\left[c_{r}^{\dagger}(\boldsymbol{p}) c_{s}\left(\boldsymbol{p}^{\prime}\right) u_{r}^{\dagger}(\boldsymbol{p}) u_{s}\left(\boldsymbol{p}^{\prime}\right) e^{i\left(\mathbf{p}-\mathbf{p}^{\prime}\right) \cdot \mathbf{x}}\right. \\
&-d_{r}(\boldsymbol{p}) d_{s}^{\dagger}\left(\boldsymbol{p}^{\prime}\right) v_{r}^{\dagger}(\boldsymbol{p}) v_{s}\left(\boldsymbol{p}^{\prime}\right) e^{-i\left(\mathbf{p}-\mathbf{p}^{\prime}\right) \cdot \mathbf{x}} \\
&-c_{r}^{\dagger}(\boldsymbol{p}) d_{s}^{\dagger}\left(\boldsymbol{p}^{\prime}\right) u_{r}^{\dagger}(\boldsymbol{p}) v_{s}\left(\boldsymbol{p}^{\prime}\right) e^{i\left(\mathbf{p}+\mathbf{p}^{\prime}\right) \cdot \mathbf{x}}  \tag{10.30}\\
&\left.+d_{r}(\boldsymbol{p}) c_{s}\left(\boldsymbol{p}^{\prime}\right) v_{r}^{\dagger}(\boldsymbol{p}) u_{s}\left(\boldsymbol{p}^{\prime}\right) e^{-i\left(\mathbf{p}+\mathbf{p}^{\prime}\right) \cdot \mathbf{x}}\right] d^{3} x \tag{10.31}
\end{align*}
$$

For the reduction of (10.30) into (10.31), the reader may refer to the solved problem 10.1.

When $t=0$, the exponential factors in Eq. (10.30) involve only threemomenta $\boldsymbol{p}$ instead of four-momenta $\mathbf{p}$ and so the exponential factors reduce to

$$
e^{i\left(\mathbf{p}-\mathbf{p}^{\prime}\right) \cdot \mathbf{x}} \xrightarrow{t=0} e^{-i\left(\boldsymbol{p}-\boldsymbol{p}^{\prime}\right) \cdot \boldsymbol{x}} .
$$

We can also replace the discrete summation over three-momenta by an integral in the case of continuous variation of three-momenta, using the following prescription.

$$
\frac{1}{V} \sum_{p} \longrightarrow \int \frac{1}{(2 \pi)^{3}} d^{3} p
$$

Let us now find the total charge that resides in the field. For this, let us first find the four-vector current $j^{\mu}(\mathbf{x})$ of the Dirac field.

$$
\begin{align*}
j^{\mu}(\mathbf{x}) & =-i e\left(\frac{\partial \mathscr{L}}{\partial \psi_{, \mu}} \psi-\frac{\partial \mathscr{L}}{\partial \psi_{, \mu}^{\dagger}} \psi^{\dagger}\right) \\
& =e \bar{\psi}(\mathbf{x}) \gamma^{\mu} \psi(\mathbf{x}) \tag{10.32}
\end{align*}
$$

The zeroth component of the four-current $j^{\mu}$ yields the charge.

$$
\begin{align*}
Q & =\int_{V} j^{0}(\mathbf{x}) d^{3} x \\
& =e \int_{V} \bar{\psi}(\mathbf{x}) \gamma^{0} \psi(\mathbf{x}) d^{3} x=e \int_{V} \psi^{\dagger}(\mathbf{x}) \psi(\mathbf{x}) d^{3} x \\
& =e \sum_{p} \sum_{r=1}^{2}\left(c_{r}^{\dagger}(\boldsymbol{p}) c_{r}(\boldsymbol{p})+d_{r}(\boldsymbol{p}) d_{r}^{\dagger}(\boldsymbol{p})\right) \tag{10.33}
\end{align*}
$$

Expressions (10.31) and (10.33) give the total energy and the total charge of the Dirac field in terms of the annihilation and creation operators. Hitherto, no commutation relation has been assumed between the creation and annihilation operators. So, let us choose either the commutation or anticommutation relation between the annihilation and creation operators so that we get physically meaningful results.

The total energy will not be positive-definite if annihilation and creation operators commute. On the other hand, the choice of anticommutation relation between annihilation and creation operator makes the energy positive definite, although the zero point energy becomes negative. This should not cause any serious problem since we can redefine the energy scale with a convenient choice of zero point energy.

### 10.4 Quantization of the Dirac field

The choice of commutation relation makes the total charge positive definite. On the other hand, the choice of anticommutation relation makes the total charge either positive or negative. This is physically acceptable. The particles and antiparticles have opposite charges and the net charge of the field is their difference.

Let us now write down the anticommutation relations between the annihilation and creation operators.

$$
\begin{array}{ll}
\left\{c_{r}(\boldsymbol{p}), c_{s}^{\dagger}\left(\boldsymbol{p}^{\prime}\right)\right\}_{+}=\delta_{r s} \delta\left(\boldsymbol{p}-\boldsymbol{p}^{\prime}\right) ; & \left\{d_{r}(\boldsymbol{p}), d_{s}^{\dagger}\left(\boldsymbol{p}^{\prime}\right)\right\}_{+}=\delta_{r s} \delta\left(\boldsymbol{p}-\boldsymbol{p}^{\prime}\right) ; \\
\left\{c_{r}(\boldsymbol{p}), c_{s}\left(\boldsymbol{p}^{\prime}\right)\right\}_{+}=0 ; & \left\{c_{r}^{\dagger}(\boldsymbol{p}), c_{s}^{\dagger}\left(\boldsymbol{p}^{\prime}\right)\right\}_{+}=0 ;  \tag{10.34}\\
\left\{d_{r}(\boldsymbol{p}), d_{s}\left(\boldsymbol{p}^{\prime}\right)\right\}_{+}=0 ; & \left\{d_{r}^{\dagger}(\boldsymbol{p}), c_{s}^{\dagger}\left(\boldsymbol{p}^{\prime}\right)\right\}_{+}=0 .
\end{array}
$$

Using the anticommutation relations (10.34), the total energy and charge
of the field become

$$
\begin{align*}
H & =\sum_{p} \sum_{r=1}^{2} E_{p}\left(c_{r}^{\dagger}(\boldsymbol{p}) c_{r}(\boldsymbol{p})+d_{r}^{\dagger}(\boldsymbol{p}) d_{r}(\boldsymbol{p})\right)-\sum_{p} \sum_{r=1}^{2} E_{p} \\
& =\sum_{p} \sum_{r=1}^{2} E_{p}\left(N_{r}(\boldsymbol{p})+\tilde{N}_{r}(\boldsymbol{p})\right)+E_{0}  \tag{10.36}\\
Q & =e \sum_{p} \sum_{r=1}^{2}\left(N_{r}(\boldsymbol{p})-\tilde{N}_{r}(\boldsymbol{p})\right)+Q_{0} \tag{10.37}
\end{align*}
$$

where $N_{r}$ and $\tilde{N}_{r}$ denote the number operator for the particles and antiparticles respectively.

$$
\begin{equation*}
N_{r}(\boldsymbol{p})=c_{r}^{\dagger}(\boldsymbol{p}) c_{r}(\boldsymbol{p}) \quad \text { and } \quad \tilde{N}_{r}(\boldsymbol{p})=d_{r}^{\dagger}(\boldsymbol{p}) d_{r}(\boldsymbol{p}) \tag{10.38}
\end{equation*}
$$

The zero-point energy and charge are denoted by $E_{0}$ and $Q_{0}$.

$$
\begin{equation*}
E_{0}=-\sum_{p} \sum_{r=1}^{2} E_{p}=-2 \sum_{p} E_{p} ; \quad Q_{0}=\sum_{p} \sum_{r=1}^{2} e=2 e \sum_{p} 1 .( \tag{10.39}
\end{equation*}
$$

If we identify the zero-point energy and charge with the ground state or the vacuum state, then what is physically observable is the energy of a system with respect to the ground state. So, the physically measurable energy and charge of a system are given by

$$
\begin{align*}
& E=\sum_{p} \sum_{r=1}^{2} E_{p}\left(n_{r}(\boldsymbol{p})+\tilde{n}_{r}(\boldsymbol{p})\right) ;  \tag{10.40}\\
& Q=e \sum_{p} \sum_{r=1}^{2}\left(n_{r}(\boldsymbol{p})-\tilde{n}_{r}(\boldsymbol{p})\right) ; \tag{10.41}
\end{align*}
$$

where $n_{r}(\boldsymbol{p})$ and $\tilde{n}_{r}(\boldsymbol{p})$ denote the eigenvalues of the number operators $N_{r}$ and $\tilde{N}_{r}$. The anticommutation relation restricts the eigenvalues of the number operator $n_{r}(\boldsymbol{p}), \tilde{n}_{r}(\boldsymbol{p})=0,1$.

The zero point energy will be avoided if we resort to the normal ordering. The normal ordering for Fermion operators is to arrange the annihilation operators to the right of creation operators, assuming the anticommutators to be zero. In the case of Bosons, the normal ordering was done by assuming the commutators to be zero.

### 10.4.1 Anticommutator between Dirac field functions

Having established the anticommutation relations (10.34) between the annihilation and creation operators, the anticommutation relation between the Dirac field operators at equal times $\psi_{\alpha}(\boldsymbol{x})$ and $\bar{\psi}_{\beta}(\boldsymbol{x})$ can be deduced. Treating $\boldsymbol{p}$ as a continuous variable and putting $t=0$, we get the following expressions from Eqs. (10.26) and (10.27).

$$
\begin{align*}
& \psi(\boldsymbol{x}, t=0) \\
& \quad=\frac{1}{(2 \pi)^{3 / 2}} \int d^{3} p \frac{1}{2 E_{p}} \sum_{r=1}^{2}\left\{c_{r}(\boldsymbol{p}) u_{r}(\boldsymbol{p}) e^{i \boldsymbol{p} \cdot \boldsymbol{x}}+d_{r}^{\dagger}(\boldsymbol{p}) v_{r}(\boldsymbol{p}) e^{-i \boldsymbol{p} \cdot \boldsymbol{x}}\right\}  \tag{10.42}\\
& \bar{\psi}(\boldsymbol{x}, t=0) \\
& \quad=\frac{1}{(2 \pi)^{3 / 2}} \int d^{3} p \frac{1}{2 E_{p}} \sum_{r=1}^{2}\left\{c_{r}^{\dagger}(\boldsymbol{p}) \bar{u}_{r}(\boldsymbol{p}) e^{-i \boldsymbol{p} \cdot \boldsymbol{x}}+d_{r}(\boldsymbol{p}) \bar{v}_{r}(\boldsymbol{p}) e^{i \boldsymbol{p} \cdot \boldsymbol{x}}\right\}(. \tag{10.43}
\end{align*}
$$

Since $\psi$ is a four-component spinor, each component has to be treated as independent field variable ${ }^{3}$. Let us now evaluate the anticommutator $\left\{\psi_{\alpha}(\boldsymbol{x}), \bar{\psi}_{\beta}(\boldsymbol{x})\right\}_{+}$.

$$
\begin{align*}
& \left\{\psi_{\alpha}(\boldsymbol{x}), \bar{\psi}_{\beta}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+} \\
& =\frac{1}{2 \pi)^{3}} \int \frac{d^{3} p d^{3} p^{\prime}}{\sqrt{2 E_{p} 2 E_{p^{\prime}}}} \sum_{r, s}\left\{\left\{c_{r}(\boldsymbol{p}), c_{s}^{\dagger}\left(\boldsymbol{p}^{\prime}\right)\right\}_{+} u_{r}^{\alpha}(\boldsymbol{p}) \bar{u}_{s}^{\beta}\left(\boldsymbol{p}^{\prime}\right) e^{i\left(\boldsymbol{p} \cdot \boldsymbol{x}-\boldsymbol{p}^{\prime} \cdot \boldsymbol{x}^{\prime}\right)}\right. \\
& \left.\quad+\left\{d_{r}^{\alpha}(\boldsymbol{p}), d_{s}^{\beta}\left(\boldsymbol{p}^{\prime}\right)\right\}_{+} v_{r}^{\alpha}(\boldsymbol{p}) \bar{v}_{s}^{\beta}\left(\boldsymbol{p}^{\prime}\right) e^{-i\left(\boldsymbol{p} \cdot \boldsymbol{x}-\boldsymbol{p}^{\prime} \cdot \boldsymbol{x}^{\prime}\right)}\right\} \\
& =  \tag{10.44}\\
& \frac{1}{2 \pi)^{3}} \int \frac{d^{3} p}{2 E_{p}} \sum_{r}\left\{u_{r}^{\alpha}(\boldsymbol{p}) \bar{u}_{r}^{\beta}(\boldsymbol{p}) e^{i \boldsymbol{p} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)}+v_{r}^{\alpha}(\boldsymbol{p}) \bar{v}_{r}^{\beta}\left(\boldsymbol{p}^{\prime}\right) e^{-i \boldsymbol{p} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)}\right\} .
\end{align*}
$$

The last step is obtained using the anticommutation relations (10.34) for the annihilation and creation operators. The spinor indices $\alpha, \beta$ are used either as subscript or superscript depending upon convenience. Since

$$
\sum_{r} u_{r}^{\alpha}(\boldsymbol{p}) \bar{u}_{r}^{\beta}(\boldsymbol{p})=(\not p+m)_{\alpha, \beta}, \quad \sum_{r} v_{r}^{\alpha}(\boldsymbol{p}) \bar{v}_{r}^{\beta}(\boldsymbol{p})=(\not p-m)_{\alpha, \beta},
$$

[^64]we can rewrite Eq. (10.44) as
\[

$$
\begin{align*}
& \left\{\psi_{\alpha}(\boldsymbol{x}), \bar{\psi}_{\beta}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+} \\
& \quad=\frac{1}{2 \pi)^{3}} \int \frac{d^{3} p}{2 E_{p}}\left\{(\not p+m)_{\alpha \beta} e^{i \boldsymbol{p} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)}+(\not p-m)_{\alpha \beta} e^{-i \boldsymbol{p} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)}\right\} . \tag{10.45}
\end{align*}
$$
\]

Expanding

$$
\not p+m=\gamma_{0} E_{p}-\gamma \cdot \boldsymbol{p}+m, \quad \not p-m=\gamma_{0} E_{p}-\gamma \cdot \boldsymbol{p}-m,
$$

and integrating over $d^{3} p$, we find that only the $\gamma_{0}$ term contributes, whereas the other terms cancel away ${ }^{4}$. Thus we get the anticommutation relations

$$
\begin{align*}
& \left\{\psi_{\alpha}(\boldsymbol{x}), \bar{\psi}_{\beta}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+}=\left(\gamma_{0}\right)_{\alpha \beta} \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right) ;  \tag{10.46}\\
& \left\{\psi_{\alpha}(\boldsymbol{x}), \psi_{\beta}^{\dagger}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+}=I_{\alpha, \beta} \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right) \tag{10.47}
\end{align*}
$$

Since $\Pi(\boldsymbol{x})=i \psi^{\dagger}(\boldsymbol{x})$ according to Eq. (10.22), we get the anticommutation relation

$$
\begin{equation*}
\left\{\psi_{\alpha}(\boldsymbol{x}), \Pi_{\beta}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+}=i I_{\alpha, \beta} \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right) \tag{10.48}
\end{equation*}
$$

where $I$ is the unit matrix.

### 10.5 Covariant anticommutation relations

Just as we have obtained the covariant commutation relations for the scalar field operators in chapter 9 , we can deduce the covariant anticommutation relations for the Dirac operators.

Let us start with the covariant form of the Dirac field operators $\psi(\mathbf{x})$ and $\bar{\psi}(\mathbf{x})$ and their Lorentz invariant divisions into positive frequency or energy parts $\psi^{+}(\mathbf{x}), \bar{\psi}^{+}(\mathbf{x})$ and negative frequency or energy parts $\psi^{-}(\mathbf{x}), \bar{\psi}^{-}(\mathbf{x})$.

$$
\begin{equation*}
\psi(\mathbf{x})=\psi^{+}(\mathbf{x})+\psi^{-}(\mathbf{x}) ; \quad \bar{\psi}(\mathbf{x})=\bar{\psi}^{+}(\mathbf{x})+\bar{\psi}^{-}(\mathbf{x}) ; \tag{10.49}
\end{equation*}
$$

$$
\begin{aligned}
& { }^{4} \text { The following integral relations will be useful: } \\
& \qquad \begin{aligned}
\int \boldsymbol{\gamma} \cdot \boldsymbol{p} e^{i \boldsymbol{p} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)} d^{3} p & =\int \boldsymbol{\gamma} \cdot \boldsymbol{p}^{\prime} e^{i \boldsymbol{p}^{\prime} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)} d^{3} p^{\prime} \text { with } \boldsymbol{p}^{\prime}=-\boldsymbol{p} . \\
\frac{1}{(2 \pi)^{3}} \int e^{+i \boldsymbol{p} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)} d^{3} p & =\frac{1}{(2 \pi)^{3}} \int e^{-i \boldsymbol{p} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)} d^{3} p=\delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right) .
\end{aligned}
\end{aligned}
$$

with

$$
\begin{align*}
\psi^{+}(\mathbf{x}) & =\frac{1}{(2 \pi)^{3 / 2}} \int_{p_{0}>0} \frac{d^{3} p}{\sqrt{2 E_{p}}} \sum_{r=1}^{2} c_{r}(\boldsymbol{p}) u_{r}(\boldsymbol{p}) e^{-i \mathbf{p} \cdot \mathbf{x}}  \tag{10.50}\\
\psi^{-}(\mathbf{x}) & =\frac{1}{(2 \pi)^{3 / 2}} \int_{p_{0}>0} \frac{d^{3} p}{\sqrt{2 E_{p}}} \sum_{r=1}^{2} d_{r}^{\dagger}(\boldsymbol{p}) v_{r}(\boldsymbol{p}) e^{i \mathbf{p} \cdot \mathbf{x}}  \tag{10.51}\\
\bar{\psi}^{+}(\mathbf{x}) & =\frac{1}{(2 \pi)^{3 / 2}} \int_{p_{0}>0} \frac{d^{3} p}{\sqrt{2 E_{p}}} \sum_{r=1}^{2} d_{r}(\boldsymbol{p}) \bar{v}_{r}(\boldsymbol{p}) e^{-i \mathbf{p} \cdot \mathbf{x}}  \tag{10.52}\\
\bar{\psi}^{-}(\mathbf{x}) & =\frac{1}{(2 \pi)^{3 / 2}} \int_{p_{0}>0} \frac{d^{3} p}{\sqrt{2 E_{p}}} \sum_{r=1}^{2} c_{r}^{\dagger}(\boldsymbol{p}) \bar{u}_{r}(\boldsymbol{p}) e^{i \mathbf{p} \cdot \mathbf{x}} \tag{10.53}
\end{align*}
$$

These operators have some physical significance. $\psi^{+}(\mathbf{x})$ is the destruction operator for the Fermion, $\psi^{-}(\mathbf{x})$ is the creation operator for the anti-Fermion, $\bar{\psi}^{+}(\mathbf{x})$ is the destruction operator for the anti-Fermion and $\bar{\psi}^{-}(\mathbf{x})$ is the creation operator for the Fermion. The physical vacuum is characterized by the property that the destruction operators operating on the vacuum yield zero.

$$
\begin{equation*}
\psi^{+}(\mathbf{x})|0\rangle=0 ; \quad \bar{\psi}^{+}(\mathbf{x})|0\rangle=0 \tag{10.54}
\end{equation*}
$$

The anticommutation rules for the creation and destruction operators are

$$
\begin{align*}
\left\{\psi^{+}(\mathbf{x}), \psi^{-}\left(\mathbf{x}^{\prime}\right)\right\}_{+} & =\left\{\bar{\psi}^{-}(\mathbf{x}), \bar{\psi}^{+}\left(\mathbf{x}^{\prime}\right)\right\}_{+}=0  \tag{10.55}\\
\left\{\bar{\psi}^{+}(\mathbf{x}), \psi^{+}\left(\mathbf{x}^{\prime}\right)\right\}_{+} & =\left\{\bar{\psi}^{-}(\mathbf{x}), \psi^{-}\left(\mathbf{x}^{\prime}\right)\right\}_{+}=0 . \tag{10.56}
\end{align*}
$$

For the anticommutator $\left\{\psi_{\alpha}^{+}(\mathbf{x}), \bar{\psi}_{\beta}^{-}(\mathbf{x})\right\}_{+}$, we have

$$
\begin{align*}
\left\{\psi_{\alpha}^{+}(\mathbf{x}), \bar{\psi}_{\beta}^{-}\left(\mathbf{x}^{\prime}\right)\right\}_{+}= & \frac{1}{(2 \pi)^{3}} \int \frac{d^{3} p}{\sqrt{2 E_{p}}} \int \frac{d^{3} p^{\prime}}{\sqrt{2 E_{p^{\prime}}}} \sum_{r=1}^{2} \sum_{s=1}^{2}\left\{c_{r}(\boldsymbol{p}), c_{s}^{\dagger}\left(\boldsymbol{p}^{\prime}\right)\right\}_{+} \\
& \times u_{r}^{\alpha}(\boldsymbol{p}) \bar{u}_{s}^{\beta}\left(\boldsymbol{p}^{\prime}\right) e^{-i \mathbf{p} \cdot \mathbf{x}^{i} e^{\prime} \mathbf{p}^{\prime} \mathbf{x}^{\prime}} . \tag{10.57}
\end{align*}
$$

Since

$$
\begin{aligned}
\left\{c_{r}(\boldsymbol{p}), c_{s}^{\dagger}\left(\boldsymbol{p}^{\prime}\right)\right\}_{+} & =\delta_{r s} \delta\left(\boldsymbol{p}-\boldsymbol{p}^{\prime}\right) \\
\sum_{r} u_{r}(\boldsymbol{p}) \bar{u}_{s}(\boldsymbol{p}) & =\not p+m
\end{aligned}
$$

we get

$$
\begin{align*}
\left\{\psi_{\alpha}^{+}(\mathbf{x}), \bar{\psi}_{\beta}^{-}\left(\mathbf{x}^{\prime}\right)\right\}_{+} & =\frac{1}{(2 \pi)^{3}} \int \frac{d^{3} p}{2 E_{p}} \sum_{r=1}^{2} u_{r}^{\alpha}(\boldsymbol{p}) \bar{u}_{r}^{\beta}(\boldsymbol{p}) e^{-i \mathbf{p} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \\
& =\frac{1}{2(2 \pi)^{3}} \int \frac{d^{3} p}{E_{p}}(\not p+m)_{\alpha \beta} e^{-i \mathbf{p} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \\
& =\left(i \nabla_{\mathbf{x}}+m\right)_{\alpha \beta} \frac{1}{2(2 \pi)^{3}} \int \frac{d^{3} p}{E_{p}} e^{-i \mathbf{p}\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \\
& =i\left(i \nabla_{\mathbf{x}}+m\right)_{\alpha \beta} \Delta^{+}\left(\mathbf{x}-\mathbf{x}^{\prime}\right) \\
& =-i S_{\alpha \beta}^{+}\left(\mathbf{x}-\mathbf{x}^{\prime}\right) \tag{10.58}
\end{align*}
$$

where

$$
\begin{align*}
\Delta^{+}\left(\mathbf{x}-\mathbf{x}^{\prime}\right) & =\frac{-i}{2(2 \pi)^{3}} \int_{p_{0}>0} \frac{d^{3} p}{p_{0}} e^{-i \mathbf{p} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}  \tag{10.59}\\
S_{\alpha \beta}^{+}\left(\mathbf{x}-\mathbf{x}^{\prime}\right) & =-\left(i \nabla_{\mathbf{x}}+m\right)_{\alpha \beta} \Delta^{+}\left(\mathbf{x}-\mathbf{x}^{\prime}\right) \tag{10.60}
\end{align*}
$$

In a similar way, we can deduce the anticommutation relation for the antiparticle operators.

$$
\begin{align*}
\left\{\psi_{\alpha}^{-}(\mathbf{x}), \bar{\psi}_{\beta}^{+}\left(\mathbf{x}^{\prime}\right)\right\}_{+}= & \frac{1}{(2 \pi)^{3}} \int \frac{d^{3} p}{\sqrt{2 E_{p}}} \int \frac{d^{3} p^{\prime}}{\sqrt{2 E_{p^{\prime}}}} \sum_{r=1}^{2} \sum_{s=1}^{2}\left\{d_{r}^{\dagger}(\boldsymbol{p}), d_{s}\left(\boldsymbol{p}^{\prime}\right)\right\}_{+} \\
& \times v_{r}^{\alpha}(\boldsymbol{p}) \bar{v}_{s}^{\beta}\left(\boldsymbol{p}^{\prime}\right) e^{i \mathbf{p} \cdot \mathbf{x}_{e^{-i} \mathbf{p}^{\prime} \cdot \mathbf{x}^{\prime}} .} \tag{10.61}
\end{align*}
$$

Since

$$
\begin{aligned}
\left\{d_{r}^{\dagger}(\boldsymbol{p}), d_{s}\left(\boldsymbol{p}^{\prime}\right)\right\}_{+} & =\delta_{r s} \delta\left(\boldsymbol{p}-\boldsymbol{p}^{\prime}\right) \\
\sum_{r} v_{r}(\boldsymbol{p}) \bar{v}_{s}(\boldsymbol{p}) & =\not p-m
\end{aligned}
$$

we get

$$
\begin{align*}
\left\{\psi_{\alpha}^{-}(\mathbf{x}), \bar{\psi}_{\beta}^{+}\left(\mathbf{x}^{\prime}\right)\right\}_{+} & =\frac{1}{(2 \pi)^{3}} \int \frac{d^{3} p}{2 E_{p}} \sum_{r=1}^{2} v_{r}^{\alpha}(\boldsymbol{p}) \bar{v}_{r}^{\beta}(\boldsymbol{p}) e^{i \mathbf{p} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \\
& =\frac{1}{2(2 \pi)^{3}} \int \frac{d^{3} p}{E_{p}}(\not p-m)_{\alpha \beta} e^{i \mathbf{p} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \\
& =\left(i \nabla_{\mathbf{x}}-m\right)_{\alpha \beta} \frac{1}{2(2 \pi)^{3}} \int \frac{d^{3} p}{E_{p}} e^{i \mathbf{p}\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \\
& =-i\left(i \nabla_{\mathbf{x}}-m\right)_{\alpha \beta} \Delta^{-}\left(\mathbf{x}-\mathbf{x}^{\prime}\right) \\
& =-i S_{\alpha \beta}^{-}\left(\mathbf{x}-\mathbf{x}^{\prime}\right) \tag{10.62}
\end{align*}
$$

where

$$
\begin{align*}
\Delta^{-}\left(\mathbf{x}-\mathbf{x}^{\prime}\right) & =\frac{i}{2(2 \pi)^{3}} \int_{p_{0}>0} \frac{d^{3} p}{p_{0}} e^{i \mathbf{p} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}  \tag{10.63}\\
S_{\alpha \beta}^{-}\left(\mathbf{x}-\mathbf{x}^{\prime}\right) & =\left(i \nabla_{\mathbf{x}}-m\right)_{\alpha \beta} \Delta^{-}\left(\mathbf{x}-\mathbf{x}^{\prime}\right) . \tag{10.64}
\end{align*}
$$

Finally, let us find the anticommutation relation between $\psi_{\alpha}(\mathbf{x})$ and $\bar{\psi}_{\beta}\left(\mathrm{x}^{\prime}\right)$.

$$
\begin{align*}
\left\{\psi_{\alpha}(\mathbf{x}), \bar{\psi}_{\beta}\left(\mathbf{x}^{\prime}\right)\right\}_{+}= & \left\{\left(\psi_{\alpha}^{+}(\mathbf{x})+\psi_{\alpha}^{-}(\mathbf{x})\right),\left(\bar{\psi}_{\beta}^{+}\left(\mathbf{x}^{\prime}\right)+\bar{\psi}_{\beta}^{-}\left(\mathbf{x}^{\prime}\right)\right)\right\}_{+} \\
= & \left\{\psi_{\alpha}^{+}(\mathbf{x}), \bar{\psi}_{\beta}^{+}\left(\mathbf{x}^{\prime}\right)\right\}_{+}+\left\{\psi_{\alpha}^{+}(\mathbf{x}), \bar{\psi}_{\beta}^{-}\left(\mathbf{x}^{\prime}\right)\right\}_{+} \\
& +\left\{\psi_{\alpha}^{-}(\mathbf{x}), \bar{\psi}_{\beta}^{+}\left(\mathbf{x}^{\prime}\right)\right\}_{+}+\left\{\psi_{\alpha}^{-}(\mathbf{x}), \bar{\psi}_{\beta}^{-}\left(\mathbf{x}^{\prime}\right)\right\}_{+} \tag{10.65}
\end{align*}
$$

Two of the above anticommutator brackets vanish.

$$
\left\{\psi_{\alpha}^{+}(\mathbf{x}), \bar{\psi}_{\beta}^{+}\left(\mathbf{x}^{\prime}\right)\right\}_{+}=\left\{\psi_{\alpha}^{-}(\mathbf{x}), \bar{\psi}_{\beta}^{-}\left(\mathbf{x}^{\prime}\right)\right\}_{+}=0 .
$$

Using the results (10.58) and (10.62) for the other two, we get

$$
\begin{equation*}
\left\{\psi_{\alpha}(\mathbf{x}), \bar{\psi}_{\beta}\left(\mathbf{x}^{\prime}\right)\right\}_{+}=-i\left[S_{\alpha \beta}^{+}\left(\mathbf{x}-\mathbf{x}^{\prime}\right)+S_{\alpha \beta}^{-}\left(\mathbf{x}-\mathbf{x}^{\prime}\right)\right] . \tag{10.66}
\end{equation*}
$$

## Review Questions

10.1 Write down the Lagrangian density for the Dirac field and obtain the field equations. Deduce expressions for the momentum density and the Hamiltonian density for the Dirac field.
10.2 Write down the Dirac field functions in terms of a complete set of single particle Dirac wave functions and arrive at an expression for the energy of an assembly of Dirac particles.
10.3 Deduce expressions for energy and charge of an assembly of Dirac particles in terms of creation and annihilation operators and show they yield meaningful results only when anticommutation relations are assumed for the annihilation and creation operators.
10.4 Given the anticommutation relations for the annihilation and creation operators for the Dirac particle

$$
\left\{c_{r}(\boldsymbol{p}), c_{s}^{\dagger}\left(\boldsymbol{p}^{\prime}\right)\right\}_{+}=\delta_{r s} \delta\left(\boldsymbol{p}-\boldsymbol{p}^{\prime}\right), \quad\left\{d_{r}(\boldsymbol{p}), d_{s}^{\dagger}\left(\boldsymbol{p}^{\prime}\right)\right\}_{+}=\delta_{r s} \delta\left(\boldsymbol{p}-\boldsymbol{p}^{\prime}\right),
$$

deduce the following anticommutation relations for the Dirac field operators:

$$
\begin{aligned}
\left\{\psi_{\alpha}(\boldsymbol{x}), \bar{\psi}_{\beta}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+} & =\gamma_{\alpha, \delta}^{0} \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right) \\
\left\{\psi_{\alpha}(\boldsymbol{x}), \Pi_{\beta}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+} & ==i I_{\alpha \beta} \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)
\end{aligned}
$$

where $I$ is the unit matrix.
10.5 Deduce the covariant anticommutation relations between the Dirac field operators:

$$
\left\{\psi_{\alpha}^{+}(\mathbf{x}), \bar{\psi}_{\beta}^{-}(\mathbf{y})\right\}_{+}, \quad\left\{\psi_{\alpha}^{-}(\mathbf{x}), \bar{\psi}_{\beta}^{+}(\mathbf{y})\right\}_{+},
$$

where the upper sign + and - indicates the positive and negative frequency part of the Dirac field operators.

## Problems

10.1 Deduce the expression (10.31) for the energy of an assembly of Dirac particles from Eq. (10.30).
10.2 Deduce the expression (10.33) for the charge of the Dirac field in terms of the creation and annihilation operators.
10.3 Given the covariant anticommutation relations (10.58) and (10.62),

$$
\begin{aligned}
\left\{\psi_{\alpha}^{+}(\mathbf{x}), \bar{\psi}_{\beta}^{-}\left(\mathbf{x}^{\prime}\right)\right\}_{+} & =\frac{1}{(2 \pi)^{3}} \int \frac{d^{3} p}{2 E_{p}} \sum_{r=1}^{2} u_{r}^{\alpha}(\boldsymbol{p}) \bar{u}_{r}^{\beta}(\boldsymbol{p}) e^{-i \mathbf{p} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \\
& =\frac{1}{2(2 \pi)^{3}} \int \frac{d^{3} p}{E_{p}}(\not p+m)_{\alpha \beta} e^{-i \mathbf{p} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \\
\left\{\psi_{\alpha}^{-}(\mathbf{x}), \bar{\psi}_{\beta}^{+}\left(\mathbf{x}^{\prime}\right)\right\}_{+} & =\frac{1}{(2 \pi)^{3}} \int \frac{d^{3} p}{2 E_{p}} \sum_{r=1}^{2} v_{r}^{\alpha}(\boldsymbol{p}) \bar{v}_{r}^{\beta}(\boldsymbol{p}) e^{i \mathbf{p} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \\
& =\frac{1}{2(2 \pi)^{3}} \int \frac{d^{3} p}{E_{p}}(\not p-m)_{\alpha \beta} e^{i \mathbf{p} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}
\end{aligned}
$$

deduce the anticommutation relation (10.46)

$$
\left\{\psi_{\alpha}(\boldsymbol{x}), \bar{\psi}_{\beta}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+}=\left(\gamma_{0}\right)_{\alpha \beta} \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)
$$

for equal times or for $t-0$

## Solutions to Problems

10.1 Evaluate expression (10.30) at time $t=0$. Then the expression involves only three-momentum $\boldsymbol{p}$ and $\boldsymbol{p}^{\prime}$.

$$
e^{i\left(\mathbf{p}-\mathbf{p}^{\prime}\right) \cdot \mathbf{x}} \xrightarrow{t=0} e^{-i\left(\boldsymbol{p}-\boldsymbol{p}^{\prime}\right) \cdot \boldsymbol{x}} .
$$

Using the orthonormal properties (10.8) - (10.12) of Dirac spinors, we get

$$
\begin{aligned}
u_{r}^{\dagger}(\boldsymbol{p}) u_{s}\left(\boldsymbol{p}^{\prime}\right) & =2 E_{p} \delta_{r s} \delta\left(\boldsymbol{p}-\boldsymbol{p}^{\prime}\right) \\
v_{r}^{\dagger}(\boldsymbol{p}) v_{s}\left(\boldsymbol{p}^{\prime}\right) & =2 E_{p} \delta_{r s} \delta\left(\boldsymbol{p}-\boldsymbol{p}^{\prime}\right) \\
u_{r}^{\dagger}(\boldsymbol{p}) v_{s}\left(\boldsymbol{p}^{\prime}\right) & =v_{r}^{\dagger}(\boldsymbol{p}) u_{s}\left(\boldsymbol{p}^{\prime}\right)=0
\end{aligned}
$$

Substituting them in expression (10.30) is equivalent to putting $r=s$ and $\boldsymbol{p}=\boldsymbol{p}^{\prime}$.

$$
\begin{aligned}
H & =\frac{1}{V} \sum_{p} \sum_{r=1}^{2} E_{p}\left[c_{r}^{\dagger}(\boldsymbol{p}) c_{r}(\boldsymbol{p})-d_{r}(\boldsymbol{p}) d_{r}^{\dagger}(\boldsymbol{p})\right] \int d^{3} x \\
& =\sum_{p} \sum_{r=1}^{2} E_{p}\left[c_{r}^{\dagger}(\boldsymbol{p}) c_{r}(\boldsymbol{p})-d_{r}(\boldsymbol{p}) d_{r}^{\dagger}(\boldsymbol{p})\right]
\end{aligned}
$$

10.2 The charge of the Dirac field is given by

$$
Q=e \int_{V} \psi^{\dagger}(\mathbf{x}) \psi(\mathbf{x}) d^{3} x
$$

where $\psi^{\dagger}(\mathbf{x}) \psi(\mathbf{x})$ is taken at the same space-time point $\mathbf{x}$. The Fourier expansions of $\psi^{\dagger}(\mathbf{x})$ and $\psi(\mathbf{x})$ are given by Eqs. (10.28) and (10.26). They are reproduced here for the sake of convenience.

$$
\begin{aligned}
\psi^{\dagger}(\mathbf{x}) & =\frac{1}{\sqrt{V}} \sum_{p} \sqrt{\frac{1}{2 E_{p}}} \sum_{r=1}^{2}\left[c_{r}^{\dagger}(\boldsymbol{p}) u_{r}^{\dagger}(\boldsymbol{p}) e^{i \mathbf{p} \cdot \mathbf{x}}+d_{r}(\boldsymbol{p}) v_{r}^{\dagger}(\boldsymbol{p}) e^{-i \mathbf{p} \cdot \mathbf{x}}\right] \\
\psi(\mathbf{x}) & =\frac{1}{\sqrt{V}} \sum_{p^{\prime}} \sqrt{\frac{1}{2 E_{p^{\prime}}}} \sum_{s=1}^{2}\left[c_{s}\left(\boldsymbol{p}^{\prime}\right) u_{s}\left(\boldsymbol{p}^{\prime}\right) e^{-i \mathbf{p}^{\prime} \cdot \mathbf{x}}+d_{s}^{\dagger}\left(\boldsymbol{p}^{\prime}\right) v_{s}\left(\boldsymbol{p}^{\prime}\right) e^{i \mathbf{p}^{\prime} \cdot \mathbf{x}}\right]
\end{aligned}
$$

In taking the product of $\psi^{\dagger}(\mathbf{x}) \psi(\mathbf{x})$, the cross terms involving the spinors $u_{r}^{\dagger}(\boldsymbol{p}) v_{s}\left(\boldsymbol{p}^{\prime}\right)$ will vanish since they are orthogonal. Retaining only the contributing terms, we get

$$
\begin{aligned}
\psi^{\dagger}(\mathbf{x}) \psi(\mathbf{x})= & \frac{1}{V} \sum_{\boldsymbol{p}, \boldsymbol{p}^{\prime}}\left(\frac{1}{2 E_{p} 2 E_{p^{\prime}}}\right)^{\frac{1}{2}} \sum_{r, s}\left[c_{r}^{\dagger}(\boldsymbol{p}) c_{s}\left(\boldsymbol{p}^{\prime}\right) u_{r}^{\dagger}(\boldsymbol{p}) u_{s}\left(\boldsymbol{p}^{\prime}\right) e^{i\left(\mathbf{p}-\mathbf{p}^{\prime}\right) \cdot \mathbf{x}}\right. \\
& \left.+d_{r}(\boldsymbol{p}) d_{s}^{\dagger}\left(\boldsymbol{p}^{\prime}\right) v_{r}^{\dagger}(\boldsymbol{p}) v_{s}\left(\boldsymbol{p}^{\prime}\right) e^{-i\left(\mathbf{p}-\mathbf{p}^{\prime}\right) \cdot \mathbf{x}}\right]
\end{aligned}
$$

Since

$$
u_{r}^{\dagger}(\boldsymbol{p}) u_{s}\left(\boldsymbol{p}^{\prime}\right)=2 E_{p} \delta_{r s} \delta\left(\boldsymbol{p}-\boldsymbol{p}^{\prime}\right) \text { and } v_{r}^{\dagger}(\boldsymbol{p}) v_{s}\left(\boldsymbol{p}^{\prime}\right)=2 E_{p} \delta_{r s} \delta\left(\boldsymbol{p}-\boldsymbol{p}^{\prime}\right)
$$

we get

$$
\psi^{\dagger}(\mathbf{x}) \psi(\mathbf{x})=\frac{1}{V} \sum_{\boldsymbol{p}} \sum_{r}\left[c_{r}^{\dagger}(\boldsymbol{p}) c_{s}(\boldsymbol{p})+d_{r}(\boldsymbol{p}) d_{r}^{\dagger}\left(v p^{\prime}\right)\right]
$$

Since $\int d^{3} x=V$ cancels with $1 / V$, we get the simplified expression for $Q$.

$$
Q=e \sum_{\boldsymbol{p}} \sum_{r}\left[c_{r}^{\dagger}(\boldsymbol{p}) c_{s}(\boldsymbol{p})+d_{r}(\boldsymbol{p}) d_{r}^{\dagger}\left(v p^{\prime}\right)\right.
$$

10.3 The anticommutator $\left\{\psi_{\alpha}(\boldsymbol{x}), \bar{\psi}_{\beta}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+}$can be expanded in term of positive and negative frequency operators.

$$
\begin{aligned}
\left\{\psi_{\alpha}(\boldsymbol{x}), \bar{\psi}_{\beta}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+}= & \left\{\left(\psi_{\alpha}^{+}(\boldsymbol{x})+\psi_{\alpha}^{-}(\boldsymbol{x})\right),\left(\bar{\psi}_{\beta}^{+}\left(\boldsymbol{x}^{\prime}\right)+\bar{\psi}_{\beta}^{-}\left(\boldsymbol{x}^{\prime}\right)\right)\right\}_{+} \\
= & \left\{\psi_{\alpha}^{+}(\boldsymbol{x}), \bar{\psi}_{\beta}^{+}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+}+\left\{\psi_{\alpha}^{+}(\boldsymbol{x}), \bar{\psi}_{\beta}^{-}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+} \\
& +\left\{\psi_{\alpha}^{-}(\boldsymbol{x}), \bar{\psi}_{\beta}^{+}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+}+\left\{\psi_{\alpha}^{-}(\boldsymbol{x}), \bar{\psi}_{\beta}^{-}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+} \\
= & \left\{\psi_{\alpha}^{+}(\boldsymbol{x}), \bar{\psi}_{\beta}^{-}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+}+\left\{\psi_{\alpha}^{-}(\boldsymbol{x}), \bar{\psi}_{\beta}^{+}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+}
\end{aligned}
$$

since

$$
\left\{\psi_{\alpha}^{+}(\boldsymbol{x}), \bar{\psi}_{\beta}^{+}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+}=\left\{\psi_{\alpha}^{-}(\boldsymbol{x}), \bar{\psi}_{\beta}^{-}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+}=0
$$

It is given that

$$
\begin{aligned}
\left\{\psi_{\alpha}^{+}(\mathbf{x}), \bar{\psi}_{\beta}^{-}\left(\mathbf{x}^{\prime}\right)\right\}_{+} & =\frac{1}{2(2 \pi)^{3}} \int \frac{d^{3} p}{E_{p}}(\not p+m)_{\alpha \beta} e^{-i \mathbf{p} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \\
\left\{\psi_{\alpha}^{-}(\mathbf{x}), \bar{\psi}_{\beta}^{+}\left(\mathbf{x}^{\prime}\right)\right\}_{+} & =\frac{1}{2(2 \pi)^{3}} \int \frac{d^{3} p}{E_{p}}(\not p-m)_{\alpha \beta} e^{i \mathbf{p} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)}
\end{aligned}
$$

For $t=0$,

$$
e^{-i \mathbf{p} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \rightarrow e^{i \boldsymbol{p} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)} ; \quad e^{i \mathbf{p} \cdot\left(\mathbf{x}-\mathbf{x}^{\prime}\right)} \rightarrow e^{-i \boldsymbol{p} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)}
$$

Please note that

$$
\not p=\gamma_{0} E-\gamma \cdot \boldsymbol{p}
$$

Using the following integral identities

$$
\begin{aligned}
\frac{1}{(2 \pi)^{3}} \int d^{3} p e^{i \boldsymbol{p} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)} & =\frac{1}{(2 \pi)^{3}} \int d^{3} p e^{-i \boldsymbol{p} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)}=\delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right) \\
\frac{1}{(2 \pi)^{3}} \int d^{3} p(\boldsymbol{\gamma} \cdot \boldsymbol{p}) e^{i \boldsymbol{p} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)} & =\frac{1}{(2 \pi)^{3}} \int d^{3} p(-\boldsymbol{\gamma} \cdot \boldsymbol{p}) e^{-i \boldsymbol{p} \cdot\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)}
\end{aligned}
$$

we find that the terms involving $\gamma \cdot \boldsymbol{p}$ and $m$ mutually cancel away and only the term involving $\gamma_{0}$ contributes. Thus we finally get

$$
\begin{aligned}
\left\{\psi_{\alpha}(\boldsymbol{x}), \bar{\psi}_{\beta}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+} & =\left\{\psi_{\alpha}^{+}(\boldsymbol{x}), \bar{\psi}_{\beta}^{-}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+}+\left\{\psi_{\alpha}^{-}(\boldsymbol{x}), \bar{\psi}_{\beta}^{+}\left(\boldsymbol{x}^{\prime}\right)\right\}_{+} \\
& =\left(\gamma_{0}\right)_{\alpha \beta} \delta\left(\boldsymbol{x}-\boldsymbol{x}^{\prime}\right)
\end{aligned}
$$

For simplicity, we have chosen $t=0$ but this relation holds for any equal times $t=t^{\prime}$.

## Chapter 11

## The Electromagnetic Field

The electromagnetic field is described by a set of Maxwell's equations which specify the electrical and magnetic field strengths $\boldsymbol{E}$ and $\boldsymbol{B}$ at each space point and their variation with time. Alternatively, it can be described in terms of scalar and vector potentials $\phi$ and $\boldsymbol{A}$. Since $\phi$ and $\boldsymbol{A}$ together yield 4 components, it is possible to represent them as a fourvector $A^{\mu}$. The latter approach yields a field equation for a particle of zero rest mass similar to the familiar Klein-Gordon (K-G) equation. Earlier we have studied the K-G equation for a scalar field that corresponds to spin-zero particle. Here we consider a K-G equation for zero rest mass but with a four-vector function that corresponds to spin-1 particle.

$$
\square A^{\mu}(\mathbf{x})=0, \quad \mu=0,1,2,3 .
$$

This is the relativistic field equation for the photon in free space.
The electrical and magnetic field strengths $\boldsymbol{E}$ and $\boldsymbol{B}$ together form 6 components and it is shown that they correspond to six independent components of an antisymmetric tensor $F^{\mu \nu}$ of rank 2 formed with $A^{\mu}$. The scalar and vector potentials $\phi(\boldsymbol{x}, t)$ and $A(\boldsymbol{x}, t)$ are not observable quantities but the electrical and magnetic field strengths $\boldsymbol{E}(\boldsymbol{x}, t)$ and $\boldsymbol{B}(\boldsymbol{x}, t)$ that are derived from the potentials are measurable quantities. However it is found that the scalar and vector potentials are not unique but there exists more than one set of potentials that yield the same $\boldsymbol{E}$ and $\boldsymbol{B}$. This is often referred to as the invariance of the electromagnetic field under gauge transformation.

### 11.1 Maxwell's Equations

First let us write down Maxwell's equations for the electromagnetic field in unrationalized Gaussian (CGS) units, using the two different notations that are in vogue for differential operators in vector algebra.

$$
\begin{array}{rll}
\operatorname{div} \boldsymbol{B}=0 & \text { or } & \boldsymbol{\nabla} \cdot \boldsymbol{B}=0 . \\
\operatorname{curl} \boldsymbol{E}+\frac{1}{c} \frac{\partial \boldsymbol{B}}{\partial t}=0 & \text { or } & \boldsymbol{\nabla} \times \boldsymbol{E}+\frac{1}{c} \frac{\partial \boldsymbol{B}}{\partial t}=0 . \\
\operatorname{curl} \boldsymbol{B}-\frac{1}{c} \frac{\partial \boldsymbol{E}}{\partial t}=\frac{4 \pi \boldsymbol{j}}{c} & \text { or } & \boldsymbol{\nabla} \times \boldsymbol{B}-\frac{1}{c} \frac{\partial \boldsymbol{E}}{\partial t}=\frac{4 \pi \boldsymbol{j}}{c} . \\
\operatorname{div} \boldsymbol{E}=4 \pi \rho & \text { or } & \boldsymbol{\nabla} \cdot \boldsymbol{E}=4 \pi \rho . \tag{11.4}
\end{array}
$$

From Eq. (11.1), it can be inferred that $\boldsymbol{B}$ can be expressed as the curl of a vector $\boldsymbol{A}$ since div curl $\boldsymbol{A}=0$, according to the vector algebra. Therefore

$$
\begin{equation*}
\boldsymbol{B}=\operatorname{curl} \boldsymbol{A}=\boldsymbol{\nabla} \times \boldsymbol{A}, \tag{11.5}
\end{equation*}
$$

where $\boldsymbol{A}$ is known as the vector potential of the electromagnetic field. Expressing $\boldsymbol{B}$ in terms of $\boldsymbol{A}$ and using the commutativity of the differential operators $\boldsymbol{\nabla}$ and $\frac{\partial}{\partial t}$, Eq. (11.2) can be rewritten as

$$
\begin{equation*}
\boldsymbol{\nabla} \times\left(\boldsymbol{E}+\frac{1}{c} \frac{\partial \boldsymbol{A}}{\partial t}\right)=0 \tag{11.6}
\end{equation*}
$$

Since curl grad $\phi=0$, the quantity within bracket in Eq. (11.6) can be written as the gradient of a scalar function $\phi$.

$$
\begin{equation*}
-\boldsymbol{\nabla} \phi=\boldsymbol{E}+\frac{1}{c} \frac{\partial \boldsymbol{A}}{\partial t}, \tag{11.7}
\end{equation*}
$$

choosing the negative sign for convenience. Thus

$$
\begin{equation*}
\boldsymbol{E}=-\boldsymbol{\nabla} \phi-\frac{1}{c} \frac{\partial \boldsymbol{A}}{\partial t} . \tag{11.8}
\end{equation*}
$$

By this procedure, we have expressed the fields $\boldsymbol{B}$ and $\boldsymbol{E}$ in terms of $\boldsymbol{A}$ and $\phi$, known as the vector and the scalar potential describing the electromagnetic field.

Now let us express Eqs. (11.3) and (11.4) in terms of $\boldsymbol{A}$ and $\phi$. Eq. (11.3) can be rewritten as (using Eqs. (11.5) and (11.8)).

$$
\begin{equation*}
\operatorname{curl} \operatorname{curl} \boldsymbol{A}-\frac{1}{c} \frac{\partial}{\partial t}\left(-\boldsymbol{\nabla} \phi-\frac{1}{c} \frac{\partial \boldsymbol{A}}{\partial t}\right)=\frac{4 \pi \boldsymbol{j}}{c} . \tag{11.9}
\end{equation*}
$$

Using the identity

$$
\begin{equation*}
\text { curl curl } \boldsymbol{A}=\operatorname{grad} \operatorname{div} \boldsymbol{A}-\boldsymbol{\nabla}^{2} \boldsymbol{A}, \tag{11.10}
\end{equation*}
$$

in vector algebra, Eq. (11.3) becomes

$$
\begin{align*}
\boldsymbol{\nabla}(\boldsymbol{\nabla} \cdot \boldsymbol{A})-\boldsymbol{\nabla}^{2} \boldsymbol{A}-\frac{1}{c} \frac{\partial}{\partial t}\left(-\boldsymbol{\nabla} \phi-\frac{1}{c} \frac{\partial \boldsymbol{A}}{\partial t}\right) & =\frac{4 \pi \boldsymbol{j}}{c} \\
\boldsymbol{\nabla}\left(\boldsymbol{\nabla} \cdot \boldsymbol{A}+\frac{1}{c} \frac{\partial \phi}{\partial t}\right)-\boldsymbol{\nabla}^{2} \boldsymbol{A}+\frac{1}{c^{2}} \frac{\partial^{2} \boldsymbol{A}}{\partial t^{2}} & =\frac{4 \pi \boldsymbol{j}}{c} . \tag{11.11}
\end{align*}
$$

In the above Equation, put

$$
\begin{equation*}
\boldsymbol{\nabla} \cdot \boldsymbol{A}+\frac{1}{c} \frac{\partial \phi}{\partial t}=0 . \tag{11.12}
\end{equation*}
$$

Then Eq. (11.11) simplifies to

$$
\begin{equation*}
\frac{1}{c^{2}} \frac{\partial^{2} \boldsymbol{A}}{\partial t^{2}}-\nabla^{2} \boldsymbol{A}=\frac{4 \pi \boldsymbol{j}}{c} \tag{11.13}
\end{equation*}
$$

Equation (11.13) is obtained under the condition (11.12), which is known as the Lorentz condition.

Substituting Eq. (11.8) for $\boldsymbol{E}$ in Eq. (11.4) and substituting the Lorentz condition

$$
\begin{equation*}
\boldsymbol{\nabla} \cdot \boldsymbol{A}=-\frac{1}{c} \frac{\partial \phi}{\partial t} \tag{11.14}
\end{equation*}
$$

Eq. (11.4) reduces to

$$
\begin{equation*}
\frac{1}{c^{2}} \frac{\partial^{2} \phi}{\partial t^{2}}-\nabla^{2} \phi=4 \pi \rho \tag{11.15}
\end{equation*}
$$

Equations (11.13) and (11.15) can be combined and written in a relativistically covariant form

$$
\begin{equation*}
\left(\frac{1}{c^{2}} \frac{\partial^{2}}{\partial t^{2}}-\nabla^{2}\right) A^{\mu}=\frac{4 \pi j^{\mu}}{c} ; \quad(\mu=0,1,2,3) \tag{11.16}
\end{equation*}
$$

using the four-vector notation. A four-vector is defined as a set of four quantities that transform in the same way as $c t, x, y, z$ when one goes from one inertial frame of reference to another (under Lorentz transformation).

$$
\begin{array}{lll}
x^{\mu} \quad\left(=x^{0}, x^{1}, x^{2}, x^{3}\right) ; & x^{0}=c t ; \\
A^{\mu} & \left(=A^{0}, A^{1}, A^{2}, A^{3}\right) ; & A^{0}=\phi ; \\
j^{\mu} & \left(=j^{0}, j^{1}, j^{2}, j^{3}\right) ; & j^{0}=c \rho .  \tag{11.17}\\
p^{\mu} & \left(=p^{0}, p^{1}, p^{2}, p^{3}\right) ; & p^{0}=E / c .
\end{array}
$$

Defining the D'Alembertian operator

$$
\nabla=\frac{1}{c^{2}} \frac{\partial^{2}}{\partial t^{2}}-\nabla^{2}
$$

Eq. (11.16) becomes

$$
\begin{equation*}
\square A^{\mu}=\frac{4 \pi j^{\mu}}{c} . \tag{11.18}
\end{equation*}
$$

Equation (11.18) can be considered as the relativistic wave equation for the photon but the $A^{\mu}$ s are subject to the Lorentz restriction (11.14).

$$
\begin{equation*}
\sum_{\mu} \frac{\partial A^{\mu}}{\partial x^{\mu}}=0 \quad \text { or } \quad \sum_{\mu} \partial_{\mu} A^{\mu}=0 \tag{11.19}
\end{equation*}
$$

using the short-hand notation $\partial_{\mu}=\partial / \partial x^{\mu}$. Potentials satisfying the Lorentz condition, Eq. (11.19) are known to be in Lorentz gauge. It is to be observed that Eq. (11.18) together with the Lorentz condition (11.19) is equivalent to Maxwell's equations. Although $A^{\mu}$ is often referred to as potential in classical physics, it really refers to the photon field in field theory.

For a free field, $\mathrm{J}^{\mu}=0$. Consequently, the electromagnetic wave equation reduces to

$$
\begin{equation*}
\square A^{\mu}=0 \quad \text { or } \quad \sum_{\nu=0}^{3} \partial^{\nu} \partial_{\nu} A^{\mu}=0 . \tag{11.20}
\end{equation*}
$$

On inspection, it can be recognized as the Klein-Gordon equation for a particle with zero mass. The solution which we sought for the KleinGordon equation was a scalar function corresponding to spin zero particle. In the present case, we have to find a solution for the four-vector function and this will correspond to the spin 1 particle.

### 11.2 Electromagnetic Field Tensor

Given the scalar and vector potentials $\phi$ and $\boldsymbol{A}$, they can be treated as the components $A^{\mu}$ of a four-component vector potential and the components of the magnetic and electric fields can be expressed as the components of a second rank antisymmetric tensor $F^{\mu \nu}$, defined by ${ }^{1}$

$$
\begin{equation*}
F^{\mu \nu}=\frac{\partial A^{\mu}}{\partial x_{\nu}}-\frac{\partial A^{\nu}}{\partial x_{\mu}}=\partial^{\nu} A^{\mu}-\partial^{\mu} A^{\nu} \tag{11.21}
\end{equation*}
$$

using the notation $\partial^{\nu}=\partial / \partial x_{\nu}$. Since this second rank tensor is antisymmetric $F^{\mu \nu}=-F^{\nu \mu}$ and the diagonal elements with $\mu=\nu$ are zero, there are only six independent components and it can be shown that they correspond to the three components of $E$ and the three components of $B$.

From Eqs. (11.5) and (11.8), we obtain

$$
\begin{align*}
& B_{x}=\frac{\partial A_{z}}{\partial y}-\frac{\partial A_{y}}{\partial z} \equiv-\partial^{2} A^{3}+\partial^{3} A^{2}=F^{23}  \tag{11.22}\\
& B_{y}=\frac{\partial A_{x}}{\partial z}-\frac{\partial A_{z}}{\partial x} \equiv-\partial^{3} A^{1}+\partial^{1} A^{3}=F^{31}  \tag{11.23}\\
& B_{z}=\frac{\partial A_{y}}{\partial x}-\frac{\partial A_{x}}{\partial y} \equiv-\partial^{1} A^{2}+\partial^{2} A^{1}=F^{12}  \tag{11.24}\\
& E_{x}=-\frac{1}{c} \frac{\partial A_{x}}{\partial t}-\frac{\partial \phi}{\partial x} \equiv-\partial^{0} A^{1}+\partial^{1} A^{0}=F^{01}  \tag{11.25}\\
& E_{y}=-\frac{1}{c} \frac{\partial A_{y}}{\partial t}-\frac{\partial \phi}{\partial y} \equiv-\partial^{0} A^{2}+\partial^{2} A^{0}=F^{02}  \tag{11.26}\\
& E_{z}=-\frac{1}{c} \frac{\partial A_{z}}{\partial t}-\frac{\partial \phi}{\partial z} \equiv-\partial^{0} A^{3}+\partial^{3} A^{0}=F^{03} \tag{11.27}
\end{align*}
$$

Collecting the above results, we can explicitly write down the elements of the field tensor in the form of a matrix, designating the rows and columns by $\mu$ and $\nu$.

$$
\begin{aligned}
& { }^{1} \text { The differential operators are defined by } \\
& \qquad \begin{aligned}
\partial_{\mu} & =\frac{\partial}{\partial x^{\mu}}=\left(\partial_{0}, \partial_{1}, \partial_{2}, \partial_{3}\right)=\left(\frac{1}{c} \frac{\partial}{\partial t}, \frac{\partial}{\partial x}, \frac{\partial}{\partial y}, \frac{\partial}{\partial z}\right) \\
& =\left(\frac{1}{c} \frac{\partial}{\partial t}, \nabla\right) . \\
\partial^{\mu} & =g^{\mu \nu} \partial_{\nu}=\left(\frac{1}{c} \frac{\partial}{\partial t},-\nabla\right) .
\end{aligned} .
\end{aligned}
$$

$$
F^{\mu \nu}=\begin{gather*}
\mu / \nu \rightarrow  \tag{11.28}\\
0 \\
1 \\
2 \\
3
\end{gather*}\left[\begin{array}{rrrr}
0 & 1 & 2 & 3 \\
-E_{x} & 0 & E_{y} & E_{z} \\
-B_{z} & -B_{y} \\
-B_{z} & B_{y} & -B_{x} & B_{x}
\end{array}\right] .
$$

In a crisp way, the field tensor can be written as ${ }^{2}$

$$
\begin{equation*}
F^{\mu \nu} \delta_{\mu \nu}=0 ; \quad F^{i j}=\epsilon_{i j k} B_{k} ; \quad F^{0 k}=-F^{k 0}=E_{k}, \tag{11.29}
\end{equation*}
$$

Maxwell's equations (11.3) and (11.4) can be combined and written in terms of the electromagnetic field tensor $F^{\mu \nu}$.

$$
\begin{equation*}
\sum_{\mu=0}^{3} \partial_{\mu} F^{\mu \nu}=\frac{4 \pi j^{\nu}}{c}, \quad \nu=0,1,2,3 \tag{11.30}
\end{equation*}
$$

For a free field, $j^{\nu}=0$. Consequently,

$$
\begin{equation*}
\sum_{\mu} \partial_{\mu} F^{\mu \nu} \equiv \sum_{\mu} \frac{\partial F^{\mu \nu}}{\partial x^{\mu}}=0, \quad \text { for a free field. } \tag{11.31}
\end{equation*}
$$

Eqs. (11.30) are equivalent to the inhomogeneous set of Maxwell's equations (11.3) and (11.4). The other set (viz., the homogeneous set) of Maxwell's equations (11.1) and (11.2) follows directly from the definition (11.21) of field tensor $F^{\mu \nu}$.

$$
\begin{equation*}
\frac{\partial F^{\mu \nu}}{\partial x_{\lambda}}+\frac{\partial F^{\nu \lambda}}{\partial x_{\mu}}+\frac{\partial F^{\lambda \mu}}{\partial x_{\nu}}=0 \tag{11.32}
\end{equation*}
$$

where $\mu, \nu, \lambda$ are any three cyclic set out of the four indices $0,1,2,3$. When $\mu, \nu, \lambda$ are chosen as $1,2,3$, Eq. (11.32) corresponds to Maxwell's equation $\boldsymbol{\nabla} \cdot \boldsymbol{B}=0$. Choosing the indices $\mu, \nu, \lambda$ as $(2,3,0),(3,0,2),(0,2,3)$ in Eq. (11.32), we obtain the other homogeneous Maxwell's equation (11.2).

[^65]Thus, we find that the very definition of the field tensors $F^{\mu \nu}$ satisfy the homogeneous set of Maxwell's equations and so we need to consider only the inhomogeneous equations (11.30) for seeking a Lagrangian formulation.

Treating $A_{\mu}$ as field quantities, the Lagrangian density for the electromagnetic field can be written as

$$
\begin{equation*}
\mathscr{L}=-\frac{F_{\lambda \rho} F^{\lambda \rho}}{16 \pi}+\frac{j_{\lambda} A^{\lambda}}{c} . \tag{11.33}
\end{equation*}
$$

The first term on the right hand side of Eq. (11.33) is for the free field and the second term denotes the interaction in terms of the four-vector current density $j_{\mu}$. It can be shown that the inhomogeneous Maxwell's Eqs. (11.30) can be obtained from the Lagrangian density (11.33) by insisting that they satisfy the Euler-Lagrange Equation,

$$
\begin{equation*}
\frac{\partial \mathscr{L}}{\partial A_{\mu}}-\partial_{\nu} \frac{\partial \mathscr{L}}{\partial\left(\partial_{\nu} A_{\mu}\right)}=0 . \tag{11.34}
\end{equation*}
$$

It follows from (11.33) that

$$
\begin{align*}
\frac{\partial \mathscr{L}}{\partial A_{\mu}} & =\frac{j_{\lambda}}{c} \delta_{\lambda \mu}  \tag{11.35}\\
\frac{\partial \mathscr{L}}{\partial\left(\frac{\partial A_{\mu}}{\partial x_{\nu}}\right)} & =-\frac{F_{\lambda \rho}}{8 \pi} \frac{\partial F_{\lambda \rho}}{\partial\left(\frac{\partial A_{\mu}}{\partial x_{\nu}}\right)} \tag{11.36}
\end{align*}
$$

Since

$$
F_{\lambda \rho}=\frac{\partial A_{\rho}}{\partial x_{\lambda}}-\frac{\partial A_{\lambda}}{\partial x_{\rho}}
$$

we obtain

$$
\begin{equation*}
\sum_{\lambda \rho} F_{\lambda \rho} \frac{\partial F_{\lambda \rho}}{\partial\left(\frac{\partial A_{\mu}}{\partial x_{\nu}}\right)}=\sum_{\lambda \rho} F_{\lambda \rho}\left(\delta_{\lambda \mu} \delta_{\rho, \nu}-\delta_{\lambda \nu} \delta_{\rho, \mu}\right)=2 F_{\mu \nu} . \tag{11.37}
\end{equation*}
$$

Substituting these results in the Euler-Lagrange Eq. (11.34), we get

$$
\begin{equation*}
\frac{j_{\mu}}{c}-\frac{1}{4 \pi} \frac{\partial F_{\mu \nu}}{\partial x_{\nu}}=0 \tag{11.38}
\end{equation*}
$$

which is identical with Eq. (11.30).
The Lagrangian density $\mathscr{L}$ as given by Eq. (11.33) is a sum of two terms, one arising from the free field and the other from the interaction
with the outside world in the form of charge and current density. The charge and current density is an explicit function of the spatial and time coordinates $x_{\mu}$ and so only in the case of free fields, there will be conserved currents in the form of stress-energy tensor $T_{\mu \nu}$, defined by

$$
\begin{align*}
T_{\mu \nu} & =\frac{\partial \mathscr{L}}{\partial\left(\frac{\partial A_{\lambda}}{\partial x_{\nu}}\right)} \frac{\partial A_{\lambda}}{\partial x_{\mu}}-\mathscr{L} \delta_{\mu \nu} \\
& =-\frac{1}{4 \pi} F_{\nu \lambda} A_{\lambda \mu}-\mathscr{L} \delta_{\mu \nu} \tag{11.39}
\end{align*}
$$

It is possible to study the electromagnetic theory without using the potentials at all but using only the field strengths. This will be more satisfactory since the field strengths are the only physically measurable quantities and not the potentials. But the use of potentials greatly simplifies the calculations and becomes a necessity when one wants to treat the interaction and quantization of the electromagnetic field.

### 11.3 Quantization of the Electromagnetic Field

The quantization of the free field ${ }^{3}$ equation (11.20) is beset with certain difficulties.

1. Equation (11.20) is equivalent to Maxwell's equations only if it is considered in conjunction with the Lorentz condition (11.19). On the other hand, the canonical quantization procedure, followed hitherto, is applicable only to the independent fields. In the present case, there are four fields but all of them are not free and they are subject to Lorentz restriction (11.19).
2. To maintain the Lorentz invariance of the four-vector field $A_{\mu}$, the introduction of the matrix tensor $g_{\mu \nu}$ imposes opposite signs to the zeroth component $A_{0}$ and the other three components $A_{1}, A_{2}, A_{3}$. This makes it difficult to treat all the four $A_{\mu}$ 's on the same footing.
First let us treat all the four $A_{\mu}$ 's as independent and Hermitian and ignore the difference in sign of $A_{0}$ and the Lorentz constraint (11.19). The difficulties 1 and 2 are overcome by a procedure, prescribed by Gupta and Bleuler ${ }^{4}$.
[^66]The plane wave solution of (11.20) is given by ${ }^{5}$

$$
\begin{equation*}
u_{\lambda}(\mathbf{x})=\frac{1}{\sqrt{V}} \frac{1}{\sqrt{2 \omega_{k}}} \epsilon^{(\lambda)}(\boldsymbol{k}) e^{-i \mathbf{k} \cdot \mathbf{x}}, \quad(\lambda=0,1,2,3) \tag{11.40}
\end{equation*}
$$

with

$$
\begin{align*}
\mathbf{k} \cdot \mathbf{x} & =\omega_{k} t-\boldsymbol{k} \cdot \boldsymbol{x} ;  \tag{11.41}\\
\mathbf{k}^{2} & =\omega_{k}^{2}-\boldsymbol{k}^{2}=0, \quad \omega_{k}=|\boldsymbol{k}|>0 \tag{11.42}
\end{align*}
$$

In Eq. (11.40), $\epsilon^{(\lambda)}(\boldsymbol{k})(\lambda=0,1,2,3)$ denote a complete set of orthonormal vectors in the Minkowski space for each value of $\boldsymbol{k}$. The field operator $A_{\mu}(\mathbf{x})$ can be expanded in terms of $u_{\lambda}(\mathbf{x})$.

$$
A_{\mu}(\mathbf{x})=\frac{1}{\sqrt{V}} \sum_{\boldsymbol{k}} \frac{1}{\sqrt{2 \omega_{k}}} \sum_{\lambda=0}^{3} \epsilon_{\mu}^{(\lambda)}(\boldsymbol{k})\left[a^{(\lambda)}(\boldsymbol{k}) e^{-i \mathbf{k} \cdot \mathbf{x}}+a^{(\lambda)^{\dagger}}(\boldsymbol{k}) e^{i \mathbf{k} \cdot \mathbf{x}}\right]_{1}
$$

In Eq. (11.43), $\epsilon_{\mu}^{(\lambda)}(\boldsymbol{k}), \lambda=0,1,2,3$ represent unit vectors in the fourdimensional space which are normalized in a Lorentz invariant way, $\epsilon^{(0)}$ being time-like and $\epsilon^{(1)}, \epsilon^{(2)}, \epsilon^{(3)}$, space-like.

$$
\begin{align*}
\epsilon^{(\lambda)} \cdot \epsilon^{(\lambda)} & =\epsilon_{\mu}^{(\lambda)} \epsilon^{\left(\lambda^{\prime}\right)}{ }_{\mu}=\epsilon_{\mu}^{(\lambda)} g^{\mu \nu} \epsilon_{\nu}^{\left(\lambda^{\prime}\right)} \\
& =\epsilon_{0}^{(\lambda)} \epsilon_{0}^{\left(\lambda^{\prime}\right)}-\epsilon_{1}^{(\lambda)} \epsilon_{1}^{\left(\lambda^{\prime}\right)}-\epsilon_{2}^{(\lambda)} \epsilon_{2}^{\left(\lambda^{\prime}\right)}-\epsilon_{3}^{(\lambda)} \epsilon_{3}^{\left(\lambda^{\prime}\right)} \\
& =g^{\lambda \lambda^{\prime}} . \tag{11.44}
\end{align*}
$$

Let us take advantage of the arbitrariness of the vector $\epsilon^{(\lambda)}$ and choose $\epsilon^{(3)}(\boldsymbol{k})$ along the direction of $\boldsymbol{k}$ and the vectors $\epsilon^{(1)}(\boldsymbol{k})$ and $\epsilon^{(2)}(\boldsymbol{k})$ perpendicular to $\boldsymbol{k}$. Further, the $x_{3}$ axis is chosen along $\boldsymbol{k}$. With this choice, $\epsilon^{(1)}$, $\epsilon^{(2)}$ denote the transverse polarization, $\epsilon^{(3)}$, the longitudinal polarization and $\epsilon^{(0)}$, the scalar or time-like polarization. The scalar and longitudinal photons are unphysical and only the transverse photons are allowed in nature. Let us define the following quantities

$$
\begin{equation*}
A_{\mu}(\boldsymbol{k})=\sum_{\lambda=0}^{3} \epsilon_{\mu}^{(\lambda)}(\boldsymbol{k}) a^{(\lambda)}(\boldsymbol{k}) ; \quad A_{\mu}^{\dagger}(\boldsymbol{k})=\sum_{\lambda=0}^{3} \epsilon_{\mu}^{(\lambda)}(\boldsymbol{k}) a^{(\lambda)}{ }^{\dagger}(\boldsymbol{k}) \tag{11.45}
\end{equation*}
$$

In terms of these quantities, the field variable $A_{\mu}(\mathbf{x})$ can be written as

$$
\begin{align*}
A_{\mu}(\mathbf{x}) & =\frac{1}{\sqrt{V}} \sum_{\boldsymbol{k}} \frac{1}{\sqrt{2 \omega_{k}}}\left[A_{\mu}(\boldsymbol{k}) e^{-i \mathbf{k} \cdot \mathbf{x}}+A_{\mu}^{\dagger}(\boldsymbol{k}) e^{i \mathbf{k} \cdot \mathbf{x}}\right]  \tag{11.46}\\
& =A_{\mu}^{(+)}(\mathbf{x})+A_{\mu}^{(-)}(\mathbf{x}) \tag{11.47}
\end{align*}
$$

[^67]where $A_{\mu}^{(+)}(\mathbf{x})$ and $A_{\mu}^{(-)}(\mathbf{x})$ represent the positive and negative frequency parts of the field operator $A_{\mu}(\mathbf{x})$.

The operators $A_{\mu}(\boldsymbol{k})$ and $A_{\nu}^{\dagger}\left(\boldsymbol{k}^{\prime}\right)$ obey the following commutation relations.

$$
\begin{equation*}
\left[A_{\mu}(\boldsymbol{k}), A_{\nu}^{\dagger}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=-g_{\mu \nu} \delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right) \tag{11.48}
\end{equation*}
$$

where $g_{\mu \nu}$ denotes the metric that we use in representing the four-vector.
The corresponding commutation relations for the annihilation and creation operators $a^{(\lambda)}$ and $a^{(\lambda)^{\dagger}}$. are

$$
\begin{align*}
& {\left[a^{(\lambda)^{\dagger}}(\boldsymbol{k}), a^{\left(\lambda^{\prime}\right)^{\dagger}}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=\left[a^{(\lambda)}(\boldsymbol{k}), a^{\left(\lambda^{\prime}\right)}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=0}  \tag{11.49}\\
& {\left[a^{(\lambda)}(\boldsymbol{k}), a^{\left(\lambda^{\prime}\right)^{\dagger}}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=-g_{\lambda \lambda^{\prime}} \delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right)} \tag{11.50}
\end{align*}
$$

Let us write more explicitly Eq. (11.50).

$$
\begin{align*}
& {\left[a^{(0)}(\boldsymbol{k}), a^{(0)^{\dagger}}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=-\delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right)}  \tag{11.51}\\
& {\left[a^{(i)}(\boldsymbol{k}), a^{(j)^{\dagger}}\left(\boldsymbol{k}^{\prime}\right)\right]_{-}=\delta_{i j} \delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right), \quad i, j=1,2,3 .} \tag{11.52}
\end{align*}
$$

The interpretation of $a$ and $a^{\dagger}$ as annihilation and creation operators for $\lambda=1,2,3$, that is for longitudinal and transverse photons, poses no problems but the relation for scalar photons will give problems because of the minus sign on the right hand side of Eq. (11.51). One problem is with respect to the norm of the scalar photon state and the other problem is with respect to the negative energy contribution by the scalar photons. Following the procedure analogous to the one adopted in the case of scalar fields, one can get the following expression for the Hamiltonian.

$$
\begin{equation*}
H=\sum_{k} \omega_{k}\left[\sum_{\lambda=1}^{3} a^{(\lambda)^{\dagger}}(\boldsymbol{k}) a^{(\lambda)}(\boldsymbol{k})-a^{(0)^{\dagger}}(\boldsymbol{k}) a^{(0)}(\boldsymbol{k})\right] . \tag{11.53}
\end{equation*}
$$

Actually, there is a subtle difference. The number density operator for the time-like photon (scalar photon) is not $a^{(0)^{\dagger}}(\boldsymbol{k}) a^{(0)}(\boldsymbol{k})$ but $-a^{(0)^{\dagger}}(\boldsymbol{k}) a^{(0)}(\boldsymbol{k})$.

### 11.4 The Gupta-Bleuler formulation

The above quantization procedure suffers from two defects.

1. It ignores the Lorentz condition.

$$
\partial_{\mu} A_{\mu}(\mathbf{x})=0
$$

2. It assumes that all the field operators $A_{\mu}(\mathbf{x})$ are Hermitian. If they are so, then the vacuum expectation value of $A_{\mu}(\boldsymbol{k}) A_{\nu}^{\dagger}\left(\boldsymbol{k}^{\prime}\right)$ is

$$
\begin{align*}
\left\langle\Phi_{0}, A_{\mu}(\boldsymbol{k}) A_{\nu}^{\dagger}\left(\boldsymbol{k}^{\prime}\right) \Phi_{0}\right\rangle & =\left\langle A_{\mu}^{\dagger}(\boldsymbol{k}) \Phi_{0}, A_{\nu}^{\dagger}\left(\boldsymbol{k}^{\prime}\right) \Phi_{0}\right\rangle \\
& =-g_{\mu \nu} \delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right), \tag{11.54}
\end{align*}
$$

which is positive for $A_{j}(\boldsymbol{k}) A_{j}^{\dagger}(\boldsymbol{k}), \quad j=1,2,3$ and negative for $A_{0}(\boldsymbol{k}) A_{0}^{\dagger}(\boldsymbol{k})$. This is an inconsistency within the theory. It is not possible to exclude the time-like photons (scalar photons) since they are required for the formulation of a relativistic invariant theory.

The above defects were rectified by Gupta and Bleuler (1) by insisting that only the expectation values of the operators need to comply with the classical relations but not the operators themselves and (2) by redefining the scalar product using a metric operator $\eta$.

The norm $|\Phi|$ of a state $\Phi$ is usually defined by

$$
\begin{equation*}
|\Phi|=\langle\Phi \mid \Phi\rangle . \tag{11.55}
\end{equation*}
$$

This is redefined using an indefinite metric operator $\eta$.

$$
\begin{equation*}
|\Phi|=\langle\Phi| \eta|\Phi\rangle . \tag{11.56}
\end{equation*}
$$

Since the norm should be real, it follows that

$$
\eta=\eta^{\dagger}
$$

The definition of the norm of a vector is somewhat relaxed. It needs no longer be positive definite. It can be positive, negative or zero. But only the first of these possibilities is physically meaningful and admits the probability interpretation.

With this modified definition of norm, the expectation value of an operator $A$ becomes

$$
\begin{equation*}
\langle A\rangle=\langle\Phi| \eta A|\Phi\rangle . \tag{11.57}
\end{equation*}
$$

Le us choose $\eta$ in such a way that it commutes with $A_{j}(\mathbf{x}), j=1,2,3$ but anticommutes with $A_{0}(\mathbf{x})$.

$$
\begin{array}{rlr}
{\left[\eta, A_{j}(\mathbf{x})\right]_{-}} & =0, & j=1,2,3  \tag{11.58}\\
\left\{\eta, A_{0}(\mathbf{x})\right\}_{+} & =0
\end{array}
$$

Also

$$
\begin{align*}
& {\left[\eta^{2}, A_{j}\right]_{-}=\eta\left[\eta, A_{j}\right]_{-}+\left[\eta, A_{j}\right]_{-} \eta=0, \quad j=1,2,3 .} \\
& {\left[\eta^{2}, A_{0}\right]_{-}=\eta\left\{\eta, A_{0}\right\}_{+}-\left\{\eta, A_{0}\right\}_{+} \eta=0 .} \tag{11.59}
\end{align*}
$$

Thus $\eta^{2}$ is a c-number that commutes with all the operators of the field and so can be chosen to yield unity.

$$
\eta^{2}=\eta \eta^{\dagger}=1
$$

### 11.4.1 The indefinite metric operator $\eta$

The metric operator $\eta$ is so chosen to obey the commutation relations (11.58) with the field operator $A_{\mu}$. From Eq. (11.46), one can obtain the commutation relations with the annihilation and creation operators ${ }^{6}$.

$$
\begin{array}{ll}
{\left[\eta, a_{j}(\boldsymbol{k})\right]} & =0, \quad j=1,2,3 .  \tag{11.60}\\
\left\{\eta, a_{0}(\boldsymbol{k})\right\} & =0 .
\end{array}
$$

Let us consider, for simplicity, a state $\Phi$ which consists of photons of momentum $\boldsymbol{k}$ only. Then, in the occupation number representation, the state $\Phi$ can be represented as

$$
\Phi=\left|n_{0}, n_{1}, n_{2}, n_{3}\right\rangle, \quad n_{\lambda}=n_{\lambda}(\boldsymbol{k}) .
$$

It is a simple exercise to obtain the matrix elements of $a_{\lambda}(\boldsymbol{k})$ and $a_{\lambda}^{\dagger}(\boldsymbol{k})$.

$$
\begin{aligned}
& \left\langle n_{0}^{\prime}, n_{1}^{\prime}, n_{2}^{\prime}, n_{3}^{\prime}\right| a_{\lambda}(\boldsymbol{k})\left|n_{0}, n_{1}, n_{2}, n_{3}\right\rangle=\sqrt{n_{\lambda}(\boldsymbol{k})} \delta_{n_{\lambda}^{\prime}, n_{\lambda}-1} \prod_{i \neq \lambda} \delta_{n_{i}, n_{i}^{\prime}} \cdot(11.61) \\
& \left\langle n_{0}, n_{1}, n_{2}, n_{3}\right| a_{\lambda}^{\dagger}(\boldsymbol{k})\left|n_{0}^{\prime}, n_{1}^{\prime}, n_{2}^{\prime}, n_{3}^{\prime}\right\rangle=\sqrt{n_{\lambda}(\boldsymbol{k})} \delta_{n_{\lambda}, n_{\lambda}^{\prime}+1} \prod_{i \neq \lambda} \delta_{n_{i}, n_{i}^{\prime}} \cdot(11.62)
\end{aligned}
$$

The commutation relations (11.60) can be explicitly written as

$$
\begin{align*}
& \left\langle n_{0}^{\prime}, n_{1}^{\prime}, n_{2}^{\prime}, n_{3}^{\prime}\right| \eta a_{j}(\boldsymbol{k})\left|n_{0}, n_{1}, n_{2}, n_{3}\right\rangle \\
& \quad=\left\langle n_{0}^{\prime}, n_{1}^{\prime}, n_{2}^{\prime}, n_{3}^{\prime}\right| a_{j}(\boldsymbol{k}) \eta\left|n_{0}, n_{1}, n_{2}, n_{3}\right\rangle, \quad j=1,2,3 .  \tag{11.63}\\
& \left\langle n_{0}^{\prime}, n_{1}^{\prime}, n_{2}^{\prime}, n_{3}^{\prime}\right| \eta a_{0}(\boldsymbol{k})\left|n_{0}, n_{1}, n_{2}, n_{3}\right\rangle \\
& \quad=-\left\langle n_{0}^{\prime}, n_{1}^{\prime}, n_{2}^{\prime}, n_{3}^{\prime}\right| a_{0}(\boldsymbol{k}) \eta\left|n_{0}, n_{1}, n_{2}, n_{3}\right\rangle . \tag{11.64}
\end{align*}
$$

[^68]From Eqs. (11.61)-(11.64), we infer that $\eta$ is a diagonal matrix whose matrix elements are explicitly given below:

$$
\begin{align*}
\left\langle\cdots n_{j} \cdots\right| \eta\left|\cdots n_{j} \cdots\right\rangle= & \left\langle\cdots n_{j}+1 \cdots\right| \eta\left|\cdots n_{j}+1 \cdots\right\rangle, \\
& j=1,2,3 .  \tag{11.65}\\
\left\langle n_{0} \cdots \cdot\right| \eta\left|n_{0} \cdots \cdots\right\rangle= & -\left\langle n_{0}+1 \cdots \cdots\right| \eta\left|n_{0}+1 \cdots \cdots\right\rangle \tag{11.66}
\end{align*}
$$

It follows that the square of the matrix element is unity.

$$
\begin{equation*}
\left.\left|\left\langle n_{0}, n_{1}, n_{2}, n_{3}\right| \eta\right| n_{0}, n_{1}, n_{2}, n_{3}\right\rangle\left.\right|^{2}=1 \tag{11.67}
\end{equation*}
$$

From the above discussion, we conclude that $\eta$ is a diagonal matrix with the matrix element

$$
\begin{equation*}
\left\langle n_{0}^{\prime}, n_{1}^{\prime}, n_{2}^{\prime}, n_{3}^{\prime}\right| \eta\left|n_{0}, n_{1}, n_{2}, n_{3}\right\rangle=(-1)^{n_{0}} \prod_{\lambda=0}^{3} \delta_{n_{\lambda}, n_{\lambda}^{\prime}} \tag{11.68}
\end{equation*}
$$

If the state $\Phi$ consists of photons of several momenta $\boldsymbol{k}$, then

$$
\begin{equation*}
\langle\Phi| \eta|\Phi\rangle=(-1)^{S_{k}} \prod_{\boldsymbol{k}} \prod_{\lambda=0}^{3} \delta_{n_{\lambda}(\boldsymbol{k}), n_{\lambda}^{\prime}(\boldsymbol{k})}, \tag{11.69}
\end{equation*}
$$

with

$$
S_{k}=\sum_{\boldsymbol{k}} n_{0}(\boldsymbol{k}) .
$$

### 11.4.2 The Lorentz condition

The Lorentz condition explains the observed transversality of the electromagnetic waves. For the quantum theory to have the correct classical limits, it is sufficient if the expectation values of the operators and not the operators themselves obey the classical equations.

The Lorentz condition is

$$
\partial^{\mu} A_{\mu}=0
$$

It is sufficient if the expectation value obeys the Lorentz condition.

$$
\begin{equation*}
\left\langle\partial^{\mu} A_{\mu}\right\rangle=\langle\Phi| \eta \partial^{\mu} A_{\mu}|\Phi\rangle=0 \tag{11.70}
\end{equation*}
$$

Expanding $A_{\mu}$ in terms of positive and negative frequency parts as shown in Eq. (11.47), we get

$$
\begin{equation*}
\langle\Phi| \eta \partial^{\mu} A_{\mu}|\Phi\rangle=\langle\Phi| \eta \partial^{\mu} A_{\mu}^{(+)}|\Phi\rangle+\langle\Phi| \eta \partial^{\mu} A_{\mu}^{(-)}|\Phi\rangle \tag{11.71}
\end{equation*}
$$

Since

$$
\begin{equation*}
\left\langle\Phi, \eta \partial^{\mu} A_{\mu}^{(-)} \Phi\right\rangle=\left\langle\eta \partial^{\mu} A_{\mu}^{(+)} \Phi, \Phi\right\rangle \tag{11.72}
\end{equation*}
$$

it is enough if the subsidiary condition

$$
\begin{equation*}
\sum_{\mu} \partial^{\mu} A_{\mu}^{(+)} \Phi=0 \tag{11.73}
\end{equation*}
$$

is obeyed. This yields the allowed states that satisfies the Lorentz condition. Let us proceed to find out the allowed states. From Eqs. (11.46) and (11.47), we get

$$
\begin{align*}
A_{\mu}^{(+)}(\mathbf{x}) & =\frac{1}{\sqrt{V}} \sum_{\boldsymbol{k}} \frac{1}{\sqrt{2 \omega_{k}}} a_{\lambda}(\boldsymbol{k}) \epsilon_{\mu}^{(\lambda)}(\boldsymbol{k}) e^{-i \mathbf{k} \cdot \mathbf{x}}  \tag{11.74}\\
\partial_{\mu} A_{\mu}^{(+)}(\mathbf{x}) & =\frac{1}{\sqrt{V}} \sum_{\boldsymbol{k}} \frac{-i k_{\mu}}{\sqrt{2 \omega_{k}}} a_{\lambda}(\boldsymbol{k}) \epsilon_{\mu}^{(\lambda)}(\boldsymbol{k}) e^{-i \mathbf{k} \cdot \mathbf{x}} \tag{11.75}
\end{align*}
$$

We have chosen the photon momentum $\boldsymbol{k}$ and the polarization direction $\epsilon^{(3)}$ along the $x_{3}$ axis. So, $\epsilon^{(1)}$ and $\epsilon^{(2)}$ denote the transverse polarizations. With this choice,

$$
k_{1}=k_{2}=0 ; \quad k_{3}=-|\boldsymbol{k}|=-\omega_{k} ; \quad k_{0}=\omega_{k} .
$$

In the coordinate system, so chosen, the subsidiary condition yields

$$
\begin{equation*}
\sum_{\mu=0}^{3} \partial^{\mu} A_{\mu}^{(+)}(\mathbf{x})=\frac{-i}{\sqrt{2 V}} \sum_{\boldsymbol{k}} \sqrt{\omega_{k}}\left(a^{(0)}(\boldsymbol{k})-a^{(3)}(\boldsymbol{k})\right) e^{-i \mathbf{k} \cdot \mathbf{x}_{(11}} \tag{11.76}
\end{equation*}
$$

The restriction is on the allowed states $\Phi$ of the field rather than a condition on the operators $A_{\mu}$, since Eq. (11.73) leads to the condition

$$
\begin{equation*}
\left(a^{(3)}(\boldsymbol{k})-a^{(0)}(\boldsymbol{k})\right) \Phi=0 . \tag{11.77}
\end{equation*}
$$

The restriction is on the combination of longitudinal and scalar photons and does not affect the transverse photons. It is possible to construct the same set of transverse photons but different combinations of longitudinal and scalar photons in accordance with Eq. (11.76). The different states so constructed are related by gauge transformations. For a free field, the gauge can be so chosen that only the transverse photons are present.

For free fields (i.e. with no charges present), the vacuum can be represented by the state $\Phi_{0}$, in which no photons of any kind are present. The
vacuum can equally be described by any state that contains no transverse photons but only an equal amount of longitudinal and scalar photons. The latter description will merely correspond to a different choice of Lorentz gauge.

The condition (11.77) makes the Hamiltonian (11.53) of the free fields positive definite, since the free fields in Lorentz gauge contain only transverse photons.

$$
H=\sum_{k} \omega_{k}\left(a_{1}^{\dagger} a_{1}+a_{2}^{\dagger} a_{2}\right) .
$$

## Review Questions

11.1 Show how Maxwell's equations for electromagnetic waves can be written in terms of scalar and vector potentials $\phi$ and $\boldsymbol{A}$. Express the electromagnetic wave equation in a relativistic covariant form by defining a four vector potential $A^{\mu}$. Show that the components of $A^{\mu}$ are not independent and they are constrained by Lorentz condition.
11.2 What are the electromagnetic field tensors? Show how Maxwell's equations can be written in terms of these field tensors.
11.3 What are the difficulties encountered in quantizing the electromagnetic field $\boldsymbol{A}^{\mu}$ and how are they overcome by the Gupta-Bleuler prescriptions?

## Problems

11.1 Given the electromagnetic field tensors $F^{\mu \nu}$ (contravariant tensor of second rank), given by Eq. (11.28), write down explicitly the matrices for $F^{\nu \mu}$ and $F_{\mu \nu}$.
11.2 The electromagnetic field is described by specifying the electrical and magnetic field strengths $\boldsymbol{E}$ and $\boldsymbol{B}$ at each point and their variation with time. It is found convenient to express them in terms of scalar and vector potentials $\phi$ and $\boldsymbol{A}$ by the following expressions:

$$
\boldsymbol{B}=\boldsymbol{\nabla} \times \boldsymbol{A} \quad \boldsymbol{E}=-\boldsymbol{\nabla} \phi-\frac{1}{c} \frac{\partial \boldsymbol{A}}{\partial t} .
$$

Show that $\phi$ and $\boldsymbol{A}$ do not form an unique set but there are more than one set that yield the same value for the physically observable quantities $\boldsymbol{B}$ and $\boldsymbol{E}$.
11.3 Given the Lagrangian density $\mathscr{L}$ of the electromagnetic field $A^{\mu}$,

$$
\mathscr{L}=-\sum_{\nu} \frac{1}{2} \partial_{\nu} A^{\mu} \partial^{\nu} A_{\mu}=-\frac{1}{2} \sum_{\nu}\left(\frac{\partial A^{\mu}}{\partial x_{\nu}}\right)^{2}
$$

obtain the momentum canonically conjugate to $A^{\mu}$ and the Hamiltonian density. Show that the field $A^{\mu}$ satisfies the field equation

$$
\square A^{\mu}=0 .
$$

## Solutions to Problems

11.1 Since $F^{\mu \nu}$ is an anti-symmetric matrix,

$$
F^{\nu \mu}=-F^{\mu \nu}
$$

$F^{\mu \nu}$ is a contravariant tensor of second rank, whereas $F_{\mu \nu}$ is a covariant tensor of second rank, obtained by lowering the indices.

$$
F_{\mu \nu}=g_{\mu \lambda} g_{\nu \rho} F^{\lambda \rho}
$$

$$
\begin{array}{lll}
F_{0 k}=g_{0 \mu} g_{k \nu} F^{\mu \nu}=-F^{0 k}, & \mu, \nu=0,1,2,3 ; \quad k=1,2,3 \\
F_{k 0}=g_{k \mu} g_{0 \nu} F^{\mu \nu}=-F^{k 0}, & \mu, \nu=0,1,2,3 ; \quad k=1,2,3 \\
F_{k l}=g_{k m} g_{l n} F^{m n}=F^{k l}, & k, l, m, n=1,2,3
\end{array}
$$

To obtain the above relations, we have used the following values of the metric tensors:

$$
g_{\mu \nu}=0, \quad \mu \neq \nu ; \quad g_{00}=1 ; \quad g_{11}=g_{22}=g_{33}=-1
$$

Using the above results, we write down explicitly the related antisymmetric tensors of second rank, given $F^{\mu \nu}$ by Eq. (11.28):

$$
\begin{aligned}
& \mu / \nu \rightarrow 0 \\
& F^{\mu \nu}= 1 \\
& 0 \\
& 1 \\
& 2 \\
& 3
\end{aligned}\left[\begin{array}{rrrr}
0 & E_{x} & E_{y} & E_{z} \\
-E_{x} & 0 & B_{z} & -B_{y} \\
-E_{y} & -B_{z} & 0 & B_{x} \\
-E_{z} & B_{y} & -B_{x} & 0
\end{array}\right] .
$$

11.2 The scalar and vector potentials $\phi$ and $\boldsymbol{A}$ are not observable quantities but only the magnetic and electrical field strengths are observable and physically measurable quantities.

$$
\boldsymbol{B}=\boldsymbol{\nabla} \times \boldsymbol{A} ; \quad \boldsymbol{E}=-\nabla \phi-\frac{1}{c} \frac{\partial \boldsymbol{A}}{\partial t}
$$

Since curl $\operatorname{div} f(\boldsymbol{x}, t)=0$, where $f(\boldsymbol{x}, t)$ is a scalar function, it is possible for two vector potentials $\boldsymbol{A}$ and $\boldsymbol{A}^{\prime}=\boldsymbol{A}-\operatorname{div} f$ to satisfy the equation for $\boldsymbol{B}$, yielding the same value for $\boldsymbol{B}$.

$$
\boldsymbol{B}=\operatorname{curl} \boldsymbol{A}=\operatorname{curl} \boldsymbol{A}^{\prime}
$$

since curl $\boldsymbol{A}^{\prime}=\operatorname{curl} \boldsymbol{A}-\operatorname{curl} \operatorname{div} f=\operatorname{curl} \boldsymbol{A}$.
For the simultaneous transformation

$$
\begin{aligned}
\boldsymbol{A} & \longrightarrow \\
\boldsymbol{A}^{\prime} & =\boldsymbol{A}-\nabla f \\
\phi & \longrightarrow
\end{aligned} \phi^{\prime}=\phi+\frac{1}{c} \frac{\partial f}{\partial t},
$$

we get

$$
\begin{aligned}
\boldsymbol{B} & =\boldsymbol{\nabla} \times \boldsymbol{A}^{\prime}=\boldsymbol{\nabla} \times \boldsymbol{A} \\
\boldsymbol{E} & =-\boldsymbol{\nabla} \phi^{\prime}-\frac{1}{c} \frac{\partial \boldsymbol{A}^{\prime}}{\partial t}=-\boldsymbol{\nabla} \phi-\frac{1}{c} \boldsymbol{\nabla} \frac{\partial f}{\partial t}-\frac{1}{c} \frac{\partial}{\partial t}(\boldsymbol{A}-\boldsymbol{\nabla} f) \\
& =-\boldsymbol{\nabla} \phi-\frac{1}{c} \frac{\partial \boldsymbol{A}}{\partial t}
\end{aligned}
$$

since $\boldsymbol{\nabla} \frac{\partial f}{\partial t}=\frac{\partial}{\partial t} \boldsymbol{\nabla} f$ because the differential operators $\boldsymbol{\nabla}$ and $\frac{\partial}{\partial t}$ commute. In the four-vector notation,

$$
A^{\mu}(\mu=0,1,2,3): \phi, \boldsymbol{A}
$$

and

$$
A^{\prime \mu}=A^{\mu}+\frac{\partial f}{\partial x^{\mu}}
$$

The last equation represents

$$
\begin{aligned}
\phi^{\prime} & =\phi+\frac{1}{c} \frac{\partial f}{\partial t} \\
\boldsymbol{A}^{\prime} & =\boldsymbol{A}-\nabla f
\end{aligned}
$$

This property is known as the invariance of electromagnetic field under gauge transformation.
11.3 Given the Lagrangian density $\mathscr{L}$ of the electromagnetic field $A^{\mu}$,

$$
\mathscr{L}=-\sum_{\nu} \frac{1}{2} \partial_{\nu} A^{\mu} \partial^{\nu} A_{\mu}=-\frac{1}{2} \sum_{\nu}\left(\frac{\partial A^{\mu}}{\partial x^{\nu}}\right)^{2}
$$

the momentum canonically conjugate to $A^{\mu}$ is given by

$$
\Pi^{\mu}=\frac{\partial \mathscr{L}}{\partial A^{\mu}} \partial t=-\frac{\partial A^{\mu}}{\partial x^{0}}=-\partial_{0} A^{\mu}
$$

The Hamiltonian density is given by

$$
\begin{aligned}
\mathscr{H} & =\Pi_{\mu} \dot{A}^{\mu}-\mathscr{L} \\
& =-\left(\frac{\partial A^{\mu}}{\partial x^{0}}\right)^{2}-\frac{1}{2} \sum_{\nu}\left(\frac{\partial A^{\mu}}{\partial x^{\nu}}\right)^{2} \\
& =-\left(\frac{\partial A^{\mu}}{\partial x^{0}}\right)^{2}+\frac{1}{2}\left(\frac{\partial A^{\mu}}{\partial x^{0}}\right)^{2}-\frac{1}{2} \sum_{k}\left(\frac{\partial A^{\mu}}{\partial x^{k}}\right)^{2} \\
& =-\frac{1}{2}\left(\frac{\partial A^{\mu}}{\partial x^{0}}\right)^{2}-\frac{1}{2} \sum_{k}\left(\frac{\partial A^{\mu}}{\partial x^{k}}\right)^{2}
\end{aligned}
$$

To deduce the field equation, we need to construct the Euler-Lagrange equation with the given Lagrangian density.

$$
\begin{equation*}
\frac{\partial \mathscr{L}}{\partial A^{\mu}}-\frac{\partial}{\partial x^{\nu}} \frac{\partial \mathscr{L}}{\partial\left(\partial_{\nu} A^{\mu}\right)}=0 \tag{11.78}
\end{equation*}
$$

With the given Lagrangian density, we have

$$
\begin{align*}
\frac{\partial \mathscr{L}}{\partial A^{\mu}} & =0  \tag{11.79}\\
\frac{\partial}{\partial x^{\nu}} \frac{\partial \mathscr{L}}{\partial\left(\partial_{\nu} A^{\mu}\right)} & =-\frac{\partial}{\partial x^{\nu}} \frac{\partial A^{\mu}}{\partial x^{\nu}}=\frac{\partial^{2} A^{\mu}}{\partial x^{\nu^{2}}} \tag{11.80}
\end{align*}
$$

Substituting (11.79) and (11.80) into Eq. (11.78), we get

$$
\begin{equation*}
\frac{\partial^{2} A^{\mu}}{\partial x^{\nu^{2}}}=0 \tag{11.81}
\end{equation*}
$$

Equation (11.81) can be written more explicitly as

$$
\frac{\partial^{2} A^{\mu}}{\partial x^{0^{2}}}-\sum_{k} \frac{\partial^{2} A^{\mu}}{\partial x^{k^{2}}}=0 \quad \text { or } \quad \square A^{\mu}=0
$$

It may be mentioned here that all the four fields $A^{\mu}, \mu=0,1,2,3$ are treated as free fields and the Lorentz condition has not been imposed.

## Chapter 12

## Interaction between Fields

Till now, we have considered only free fields but all the phenomena that we observe arise from interaction of fields. Interaction between electron and photon fields give rise to a variety of phenomena such as Compton effect, electron-positron pair creation, electron-positron pair annihilation, electron-electron scattering, Bremmstrahlung and a host of other processes, that we have discussed in Chap. 6.

It is an extremely difficult problem to solve the coupled non-linear field equations or the field equations with a source term. So, attempts have been made to solve the problem of interacting fields using perturbation theory. The Hamiltonian of the system is divided into that of the free fields plus an interaction Hamiltonian which is treated as a perturbation. This is justifiable if the interaction is sufficiently weak as in the case of electromagnetic interaction with a coupling constant $\alpha \approx 1 / 137$. This approach has been highly successful in treating the quantum electrodynamic processes.

The Lagrangian density for free fields consists of terms which are quadratic in field functions which are Fourier analyzed in terms of annihilation and creation operators but no perceptible changes will be observed. In the case of interaction between fields, one type of particles will be annihilated and either the same type of particles with different momenta and spin states or another type of particles will be created leading to observable changes in the system. The total Lagrangian $L$ and the Lagrangian density $\mathscr{L}$ can be written as a sum of free field Lagrangian $\left(L_{0}, \mathscr{L}_{0}\right)$ and
interaction Lagrangian $\left(L_{i}, \mathscr{L}_{i}\right)$.

$$
\begin{equation*}
L=L_{0}+L_{i} ; \quad \mathscr{L}=\mathscr{L}_{0}+\mathscr{L}_{i} . \tag{12.1}
\end{equation*}
$$

In a similar way, one can write the Hamiltonian $H$ and the Hamiltonian density $\mathscr{H}$ as a sum of the free-field Hamiltonian and the interacting field Hamiltonian.

$$
\begin{equation*}
H=H_{0}+H_{i} ; \quad \mathscr{H}=\mathscr{H}_{0}+\mathscr{H}_{i} . \tag{12.2}
\end{equation*}
$$

Then, it is possible to use the time-dependent perturbation formalism that we have studied in Quantum Mechanics ${ }^{1}$ for the development of perturbation theory. The free-field Hamiltonian $H_{0}$ can be considered as the unperturbed Hamiltonian and the interaction Hamiltonian $H_{i}$ can be treated as the perturbation. The interaction picture is found to be most convenient for the development of the perturbation theory.

### 12.1 The Tomonaga-Schwinger Equation

Consider the state vector $\Phi$ that obeys the Schrödinger equation.

$$
\begin{equation*}
\left(H_{0}+H_{i}\right) \Phi(t)=i \hbar \frac{\partial \Phi(t)}{\partial t} \tag{12.3}
\end{equation*}
$$

where $H_{0}$ denotes the Hamiltonian of the non-interacting fields (i.e. free fields) and $H_{i}$ denotes the interaction between fields. In the Schrödinger picture, the state vector $\Phi(t)$ is time-dependent and the operators $H_{0}, H_{i}$ are time-independent. The state vector $\Phi$, in the absence of interaction ( $H_{i}=0$ ) describes a situation in which a set of definite number of free particles with definite momentum and spin travel through space without any interaction with one another. The interaction Hamiltonian $H_{i}$ causes interaction of these particles with each other and among themselves.

Let

$$
\begin{equation*}
\Psi(t)=e^{i H_{0} t / \hbar} \Phi(t) \tag{12.4}
\end{equation*}
$$

Then differentiating (12.4), we get

$$
\begin{equation*}
i \hbar \frac{\partial \Psi(t)}{\partial t}=H_{I}(t) \Psi(t) \tag{12.5}
\end{equation*}
$$

[^69]where
\[

$$
\begin{equation*}
H_{I}(t)=e^{i H_{0} t / \hbar} H_{i} e^{-i H_{0} t / \hbar} \tag{12.6}
\end{equation*}
$$

\]

is the interaction energy operator in the interaction picture. $H_{I}(t)$ is explicitly time-dependent in contrast to the Schrödinger operator $H_{i}$ which is time-independent. In general, the operator $\mathcal{O}_{I}$ in the interaction picture is related to the Schrödinger operator $\mathcal{O}_{S}$ by the relation

$$
\begin{equation*}
\mathcal{O}_{I}=e^{i H_{0} t / \hbar} \mathcal{O}_{S} e^{-i H_{0} t / \hbar} \tag{12.7}
\end{equation*}
$$

which, on differentiation, yields the relation

$$
\begin{equation*}
i \hbar \frac{\partial \mathcal{O}_{I}(t)}{\partial t}=\left[\mathcal{O}_{I}, H_{0}\right] \tag{12.8}
\end{equation*}
$$

If we take the operator $\mathcal{O}=H_{0}$, then $H_{0}$ is the same in both the pictures, Schrödinger and interaction.

$$
\left(H_{0}\right)_{\text {int }}=\left(H_{0}\right)_{S} .
$$

The Hamiltonian $H_{0}$ contains only free fields. This assures that the field operators satisfy the free-field equations in the interaction picture also whereas the time-dependence of the state vector $\Psi(t)$ is determined by the interaction energy $H_{I}(t)$ only.

The interaction picture has distinct advantages. 1. Since the field operators satisfy free field equations, invariant commutation relations can be written down for all times. 2. Equation (12.5) can be generalized to hold in four-dimensional space-time so as to make it co-variant under Lorentz transformation to yield the Tomonaga-Schwinger equation ${ }^{2}$.

$$
\begin{equation*}
i \hbar c \frac{\delta \Psi(\boldsymbol{\sigma})}{\delta \boldsymbol{\sigma}(\mathbf{x})}=\mathscr{H}_{I}(\mathbf{x}) \Psi(\boldsymbol{\sigma}) \tag{12.9}
\end{equation*}
$$

where $\mathscr{H}_{I}(\mathbf{x})$ denotes the Hamiltonian density at space-time point $\mathbf{x}$ and $\sigma$ represents the space-like surfaces in the four-dimensional space-time.

### 12.2 The invariant perturbation theory

The differential equation (12.5) can be written as an integral equation.

$$
\begin{equation*}
\Psi(t)=\Psi\left(t_{0}\right)-\frac{i}{\hbar} \int_{t_{0}}^{t} d t^{\prime} H_{I}\left(t^{\prime}\right) \Psi\left(t^{\prime}\right) . \tag{12.10}
\end{equation*}
$$

[^70]The integral equation incorporates the initial condition that at $t=t_{0}$, the state vector $\Psi(t)=\Psi\left(t_{0}\right)$. One can define a time-development operator $U\left(t, t_{0}\right)$ such that

$$
\begin{equation*}
\Psi(t)=U\left(t, t_{0}\right) \Psi\left(t_{0}\right), \quad \Psi\left(t_{0}\right)=U^{-1}\left(t, t_{0}\right) \Psi(t) \tag{12.11}
\end{equation*}
$$

with

$$
U\left(t_{0}, t_{0}\right)=1 .
$$

The time-development operator $U\left(t, t_{0}\right)$ is a unitary operator that preserves the normalization of the state vectors $\Psi(t)$. The operator $U\left(t, t_{0}\right)$ satisfies the following differential equation

$$
\begin{equation*}
i \hbar \frac{\partial U\left(t, t_{0}\right)}{\partial t}=H_{I}(t) U\left(t, t_{0}\right), \tag{12.12}
\end{equation*}
$$

the hermitian conjugate of which can be written as

$$
\begin{equation*}
-i \hbar \frac{\partial U^{\dagger}\left(t, t_{0}\right)}{\partial t}=U^{\dagger}\left(t, t_{0}\right) H_{I}(t) \tag{12.13}
\end{equation*}
$$

since $H_{I}(t)$ is Hermitian. It can be easily verified that $U\left(t, t_{0}\right)$ is a unitary operator

$$
\begin{equation*}
U\left(t, t_{0}\right)=U^{\dagger}\left(t_{0}, t\right)=U^{-1}\left(t_{0}, t\right) \tag{12.14}
\end{equation*}
$$

and obeys the group property

$$
\begin{equation*}
U\left(t, t_{0}\right)=U\left(t, t_{1}\right) U\left(t_{1}, t_{0}\right) \tag{12.15}
\end{equation*}
$$

The differential equation (12.12) can be written as an integral equation.

$$
\begin{align*}
\int_{t_{0}}^{t} \frac{d U\left(t_{1}, t_{0}\right)}{d t_{1}} d t_{1} & =\frac{1}{i \hbar} \int_{t_{0}}^{t} H_{I}\left(t_{1}\right) U\left(t_{1}, t_{0}\right) d t_{1}, \\
U\left(t, t_{0}\right)-U\left(t_{0}, t_{0}\right) & =\frac{1}{i \hbar} \int_{t_{0}}^{t} H_{I}\left(t_{1}\right) U\left(t_{1}, t_{0}\right) d t_{1}, \\
U\left(t, t_{0}\right) & =1+\frac{1}{i \hbar} \int_{t_{0}}^{t} H_{I}\left(t_{1}\right) U\left(t_{1}, t_{0}\right) d t_{1} . \tag{12.16}
\end{align*}
$$

The last step is obtained by substituting the value $U\left(t_{0}, t_{0}\right)=1$. Equation (12.16) is the integral equation for $U\left(t, t_{0}\right)$ and by repeated iteration of $U\left(t, t_{0}\right)$ that occurs within the integral, we obtain an infinite series. Since
$H_{I}(t)$ is a small perturbation, the successive terms contribute less and less and so the series can be terminated with second or third order justifiably.

$$
\begin{align*}
U\left(t, t_{0}\right)= & 1+\frac{1}{i \hbar} \int_{t_{0}}^{t} H_{I}\left(t_{1}\right) d t_{1} \\
& +\left(\frac{1}{i \hbar}\right)^{2} \int_{t_{0}}^{t} d t_{1} \int_{t_{0}}^{t_{1}} d t_{2} H_{I}\left(t_{1}\right) H_{I}\left(t_{2}\right) U\left(t_{2}, t_{0}\right) \\
= & 1+\frac{1}{i \hbar} \int_{t_{0}}^{t} H_{I}\left(t_{1}\right) d t_{1} \\
& +\left(\frac{1}{i \hbar}\right)^{2} \int_{t_{0}}^{t} d t_{1} \int_{t_{0}}^{t_{1}} d t_{2} H_{I}\left(t_{1}\right) H_{I}\left(t_{2}\right)+\cdots . \tag{12.17}
\end{align*}
$$

The infinite series (12.17) can be written in a compact form as shown below.

$$
\begin{equation*}
U\left(t, t_{0}\right)=\sum_{n=0}^{\infty} U^{(n)}\left(t, t_{0}\right), \tag{12.18}
\end{equation*}
$$

where

$$
\begin{equation*}
U^{(0)}\left(t, t_{0}\right)=1 \tag{12.19}
\end{equation*}
$$

and

$$
\begin{equation*}
U^{(n)}\left(t, t_{0}\right)=(i \hbar)^{-n} \int_{t_{0}}^{t} d t_{1} \int_{t_{0}}^{t_{1}} d t_{2} \cdots \int_{t_{0}}^{t_{n-1}} d t_{n} H_{I}\left(t_{1}\right) H_{I}\left(t_{2}\right) \cdots H_{I}\left(t_{n}\right), \tag{12.20}
\end{equation*}
$$

with

$$
\begin{equation*}
t_{1}>t_{2}>\cdots>t_{n-1}>t_{n} . \tag{12.21}
\end{equation*}
$$

Dyson ${ }^{3}$ observed that this is essentially an integral over the whole time interval from $t_{0}$ to $t$ with the restriction that $t_{j}$ be earlier than $t_{j-1}, j \leq n$. One can remove this restriction by introducing the chronological operator $P$ which rearranges a product of time-labeled operators in a chronological time sequence with the latest time occurring first.

$$
\begin{align*}
P\left(H_{I}\left(t_{1}\right) H_{I}\left(t_{2}\right) \cdots H_{I}\left(t_{n}\right)\right)= & H_{I}\left(t_{i}\right) H_{I}\left(t_{j}\right) \cdots H_{I}\left(t_{k}\right), \\
& t_{i}>t_{j}>\cdots>t_{k} . \tag{12.22}
\end{align*}
$$

[^71]The symbol $P$ denotes Dyson's Chronological Operator. Using the symmetry of the integrand and the fact that there are $n$ ! possible permutations of the ordering of time, we obtain

$$
\begin{equation*}
U^{(n)}\left(t, t_{0}\right)=\frac{(i \hbar)^{-n}}{n!} \int_{t_{0}}^{t} d t_{1} \int_{t_{0}}^{t} d t_{2} \cdots \int_{t_{0}}^{t} d t_{n} P\left(H_{I}\left(t_{1}\right) H_{I}\left(t_{2}\right) \cdots H_{I}\left(t_{n}\right)\right) .( \tag{12.23}
\end{equation*}
$$

Thus
$U\left(t, t_{0}\right)=\sum_{n=0}^{\infty} \frac{(i \hbar)^{-n}}{n!} \int_{t_{0}}^{t} d t_{1} \int_{t_{0}}^{t} d t_{2} \cdots \int_{t_{0}}^{t} d t_{n} P\left(H_{I}\left(t_{1}\right) H_{I}\left(t_{2}\right) \cdots H_{I}\left(t_{n}\right)\right)$.
Equation (12.24) can be rewritten in a form which indicates the covariance better by expressing the interaction Hamiltonian $H_{I}(t)$ in terms of interaction Hamiltonian density $\mathscr{H}_{I}(\mathbf{x})$

$$
H_{I}(t)=\int d^{3} x \mathscr{H}_{I}(\mathbf{x})
$$

and converting the integral over time into a four-dimensional integral since $\int d t d^{3} x=\int d^{4} x / c$.

$$
\begin{align*}
U\left(t, t_{0}\right)= & \sum_{n=0}^{\infty} \frac{(i \hbar c)^{-n}}{n!} \int d^{4} x_{1} \int d^{4} x_{2} \cdots \int d^{4} x_{n} \\
& \times P\left(\mathscr{H}_{I}\left(\mathbf{x}_{1}\right) \mathscr{H}_{I}\left(\mathbf{x}_{2}\right) \cdots \mathscr{H}_{I}\left(\mathbf{x}_{n}\right)\right) \tag{12.25}
\end{align*}
$$

where $\mathbf{x}_{1}, \mathbf{x}_{2}, \cdots$ denote the space-time coordinates. Equation (12.25) forms the basis of the modern covariant perturbation theory.

### 12.3 The S-matrix

The time development operator $U\left(t, t_{0}\right)$ with $t_{0} \rightarrow-\infty$ and $t \rightarrow+\infty$ is often called the $S$-matrix or collision operator. It is also called the scattering matrix.

$$
\begin{equation*}
S=U(+\infty,-\infty) \quad \text { or } \quad \Psi(+\infty)=S \Psi(-\infty) \tag{12.26}
\end{equation*}
$$

where $\Psi(-\infty)$ and $\Psi(+\infty)$ define asymptotically the incoming and outgoing states.

We usually specify the initial and final states by the eigenfunctions of $H_{0}$ since $H_{I}$ is zero at $t= \pm \infty$. We invoke the adiabatic hypothesis to
slowly switch on the interaction $H_{I}$ at $t=-\infty$ to allow for self-interaction of the particles although they are well-separated. This process is known as "dressing up" of the particles with their virtual quanta and develop into real physical particles. Then the particles interact and make a transition to the final state. Then the interaction Hamiltonian is slowly switched off and the final state emerges as a different eigenstate of the unperturbed Hamiltonian $H_{0}$.

If $\Phi_{i}$ and $\Phi_{f}$ denote the initial and final state of the system, then

$$
\begin{align*}
S_{f i} & =\left\langle\Phi_{f}\right| S\left|\Phi_{i}\right\rangle=\left\langle\Phi_{f}\right| U(+\infty,-\infty)\left|\Phi_{i}\right\rangle \\
& =\left\langle\Phi_{f}\right| U(+\infty, 0) U(0,-\infty)\left|\Phi_{i}\right\rangle \\
& =\left\langle\Psi_{f}^{-} \mid \Psi_{i}^{+}\right\rangle \tag{12.27}
\end{align*}
$$

where $\Psi_{i}^{+}$denotes the incoming state and $\Psi_{f}^{-}$, the outgoing state.

### 12.4 Reduction of the S-matrix

### 12.4.1 Decomposition into normal products

Each term in the perturbation expansion (12.25) of the S-matrix can give rise to a variety of processes, both virtual as well as real. Dyson ${ }^{4}$ and Wick ${ }^{5}$ have shown how to express the chronological product

$$
P\left(\mathscr{H}_{I}\left(\mathbf{x}_{1}\right) \mathscr{H}_{I}\left(\mathbf{x}_{2}\right) \cdots \mathscr{H}_{I}\left(\mathbf{x}_{n}\right)\right)
$$

that occurs in the S-matrix in a form that represents explicitly all the physical processes. This is known as the decomposition of the chronological product into normal products. The normal product is defined as the product of free-particle annihilation and creation operators in which all the annihilation operators appear to the right of creation operators. Then, for the given initial and final states, there exists only one normal product that gives a non-zero matrix element between the specified initial and final states. Then the decomposition of S-matrix into normal products is to list all the matrix elements of $S$ in a representation in which all the free particle occupation numbers are diagonal. Then a Feynman diagram is simply a concise way of representing a normal product.

Let us now consider how to reduce the S-matrix into normal products, following the algebraic method of Wick. Consider a product of annihilation and creation operators. To arrange it as a normal product, all the

[^72]creation operators should be shifted to the left of annihilation operators. This can be done using the commutation relations for scalar field (Boson) operators and anticommutation relations for Dirac field (Fermion) operators, assuming that the commutators and anticommutators are zero. The effect is to include a change in sign when the anticommuting Dirac field operators are shifted from one position to another. Denoting the Scalar field operators by $\phi$ and the Dirac field operators by $\psi$ and expressing them as a sum of positive and negative frequency parts $\phi=\phi^{+}+\phi^{-}$and $\psi=\psi^{+}+\psi^{-}$, we get $^{6}$
\[

$$
\begin{align*}
N\left(\phi^{+}(\mathbf{x}) \phi^{-}(\mathbf{y})\right) & =N\left(\phi^{-}(\mathbf{y}) \phi^{+}(\mathbf{x})\right) \\
& =\phi^{-}(\mathbf{y}) \phi^{+}(\mathbf{x})  \tag{12.28}\\
N\left(\psi_{\alpha}^{+}(\mathbf{x}) \bar{\psi}_{\beta}^{-}(\mathbf{y})\right) & =-N\left(\bar{\psi}_{\beta}^{-}(\mathbf{y}) \psi_{\alpha}^{+}(\mathbf{x})\right) \\
& =-\bar{\psi}_{\beta}^{-}(\mathbf{y}) \psi_{\alpha}^{+}(\mathbf{x})  \tag{12.29}\\
N\left(\psi_{\alpha}^{+}(\mathbf{x}) \psi_{\beta}^{+}(\mathbf{y})\right) & =-N\left(\psi_{\beta}^{+}(\mathbf{y}) \psi_{\alpha}^{+}(\mathbf{x})\right) \\
& =\psi_{\alpha}^{+}(\mathbf{x}) \psi_{\beta}^{+}(\mathbf{y})=-\psi_{\beta}^{+}(\mathbf{y}) \psi_{\alpha}^{+}(\mathbf{x}) \tag{12.30}
\end{align*}
$$
\]

The normal product operation obeys the distribution law.

$$
\begin{equation*}
N\left(\phi^{+}(\mathbf{x})\left\{\phi^{+}(\mathbf{y})+\phi^{-}(\mathbf{z})\right\}\right)=N\left(\phi^{+}(\mathbf{x}) \phi^{+}(\mathbf{y})\right)+N\left(\phi^{+}(\mathbf{x}) \phi^{-}(\mathbf{z})\right) . \tag{12.31}
\end{equation*}
$$

Using the commutation and anticommutation rules, the definition of normal product and the distribution law, one can arrange the product of any two operators as normal product.

$$
\begin{align*}
\phi(\mathbf{x}) \phi(\mathbf{y})= & \left(\phi^{+}(\mathbf{x})+\phi^{-}(\mathbf{x})\right)\left(\phi^{+}(\mathbf{y})+\phi^{-}(\mathbf{y})\right) \\
= & \phi^{+}(\mathbf{x}) \phi^{+}(\mathbf{y})+\phi^{+}(\mathbf{x}) \phi^{-}(\mathbf{y})+\phi^{-}(\mathbf{x}) \phi^{+}(\mathbf{y})+\phi^{-}(\mathbf{x}) \phi^{-}(\mathbf{y}) \\
= & \phi^{+}(\mathbf{x}) \phi^{+}(\mathbf{y})+\phi^{-}(\mathbf{y}) \phi^{+}(\mathbf{x})+\left[\phi^{+}(\mathbf{x}), \phi^{-}(\mathbf{y})\right]_{-} \\
& +\phi^{-}(\mathbf{x}) \phi^{+}(\mathbf{y})+\phi^{-}(\mathbf{x}) \phi^{-}(\mathbf{y}) \\
= & N\left(\phi^{+}(\mathbf{x}) \phi^{+}(\mathbf{y})\right)+N\left(\phi^{-}(\mathbf{x}) \phi^{+}(\mathbf{y})\right)+N\left(\phi^{-}(\mathbf{x}) \phi^{+}(\mathbf{y})\right) \\
& +N\left(\phi^{-}(\mathbf{x}) \phi^{-}(\mathbf{y})\right)+\left[\phi^{+}(\mathbf{x}), \phi^{-}(\mathbf{y})\right]_{-} \\
= & N(\phi(\mathbf{x}) \phi(\mathbf{y}))+\left[\phi^{+}(\mathbf{x}), \phi^{-}(\mathbf{y})\right]_{-} \\
= & N(\phi(\mathbf{x}) \phi(\mathbf{y}))+i \hbar c \Delta^{+}(\mathbf{x}-\mathbf{y}) . \tag{12.32}
\end{align*}
$$

[^73]In deducing Eq. (12.32), we have used the relation

$$
\phi^{+}(\mathbf{x}) \phi^{-}(\mathbf{y})=\phi^{-}(\mathbf{y}) \phi^{+}(\mathbf{x})+\left[\phi^{+}(\mathbf{x}), \phi^{-}(\mathbf{y})\right]_{-}
$$

in the third step and the commutation relation

$$
\left[\phi^{+}(\mathbf{x}), \phi^{-}(\mathbf{y})\right]_{-}=i \hbar c \Delta^{+}(\mathbf{x}-\mathbf{y})
$$

in the last step.
Similarly, for fermion operators, one gets

$$
\begin{align*}
\bar{\psi}_{\alpha}(\mathbf{x}) \psi_{\beta}(\mathbf{y})= & \left(\bar{\psi}_{\alpha}^{+}(\mathbf{x})+\bar{\psi}_{\alpha}^{-}(\mathbf{x})\right)\left(\psi_{\beta}^{+}(\mathbf{y})+\psi_{\beta}^{-}(\mathbf{y})\right) \\
= & \bar{\psi}_{\alpha}^{+}(\mathbf{x}) \psi_{\beta}^{+}(\mathbf{y})+\bar{\psi}_{\alpha}^{+}(\mathbf{x}) \psi_{\beta}^{-}(\mathbf{y})+\bar{\psi}_{\alpha}^{-}(\mathbf{x}) \psi_{\beta}^{+}(\mathbf{y}) \\
& +\bar{\psi}_{\alpha}^{-}(\mathbf{x}) \psi_{\beta}^{-}(\mathbf{y}) \\
= & \bar{\psi}_{\alpha}^{+}(\mathbf{x}) \psi_{\beta}^{+}(\mathbf{y})-\psi_{\beta}^{-}(\mathbf{y}) \bar{\psi}_{\alpha}^{+}(\mathbf{x})+\left\{\bar{\psi}_{\alpha}^{+}(\mathbf{x}), \psi_{\beta}^{+}(\mathbf{y})\right\}_{+} \\
& +\bar{\psi}_{\alpha}^{-}(\mathbf{x}) \psi_{\beta}^{+}(\mathbf{y})+\bar{\psi}_{\alpha}^{-}(\mathbf{x}) \psi_{\beta}^{-}(\mathbf{y}) \\
= & N\left(\bar{\psi}_{\alpha}^{+}(\mathbf{x}) \psi_{\beta}^{+}(\mathbf{y})\right)+N\left(\bar{\psi}_{\alpha}^{+}(\mathbf{x}) \psi_{\beta}^{-}(\mathbf{y})\right) \\
& +N\left(\bar{\psi}_{\alpha}^{-}(\mathbf{x}) \psi_{\beta}^{+}(\mathbf{y})\right)+N\left(\bar{\psi}_{\alpha}^{-}(\mathbf{x}) \psi_{\beta}^{-}(\mathbf{y})\right) \\
& +\left\{\bar{\psi}_{\alpha}^{+}(\mathbf{x}), \psi_{\beta}^{+}(\mathbf{y})\right\}_{+} \\
= & N\left(\bar{\psi}_{\alpha}(\mathbf{x}) \psi_{\beta}(\mathbf{y})\right)-i S_{\beta \alpha}^{-}(\mathbf{y}-\mathbf{x}), \tag{12.33}
\end{align*}
$$

using the anticommutation relation $\left\{\bar{\psi}_{\alpha}^{+}(\mathbf{x}), \psi_{\beta}^{+}(\mathbf{y})\right\}_{+}=-i S_{\beta \alpha}^{-}(\mathbf{y}-\mathbf{x})$ for Fermion operators.

Using similar algebraic manipulations, we get

$$
\begin{align*}
\psi_{\alpha}(\mathbf{x}) \bar{\psi}_{\beta}(\mathbf{y}) & =N\left(\psi_{\alpha}(\mathbf{x}) \bar{\psi}_{\beta}(\mathbf{y})\right)-i S_{\alpha \beta}^{+}(\mathbf{x}-\mathbf{y}) ;  \tag{12.34}\\
\psi(\mathbf{x}) \psi(\mathbf{y}) & =N(\psi(\mathbf{x}) \psi(\mathbf{y}))  \tag{12.35}\\
\bar{\psi}(\mathbf{x}) \bar{\psi}(\mathbf{y}) & =N(\bar{\psi}(\mathbf{x}) \bar{\psi}(\mathbf{y})) ;  \tag{12.36}\\
\phi(\mathbf{x}) \psi(\mathbf{y}) & =N(\phi(\mathbf{x}) \psi(\mathbf{y})) \tag{12.37}
\end{align*}
$$

### 12.4.2 Wick's chronological product

Wick's chronological operator $T$ is similar to Dyson's chronological operator $P$ and arranges the time-labeled field operators $A, B, C, D \cdots$ according to the time sequence, the operators having the earlier times on the right and the operators having the later times on the left. The only difference between these two chronological operators is that Wick's chronological operator includes in its definition the sign of the permutation of

Fermion operators in making such an arrangement.

$$
\begin{equation*}
T(A B C D \cdots)=\delta_{P} P(A B C D \cdots) \tag{12.38}
\end{equation*}
$$

The whole product carries the plus or minus sign $\delta_{P}$ according to the even or odd number of permutations that the Fermion operators have to make in going from one side to the other in the process of time-ordering.

We wish to reiterate that the chronological operator $T$ arranges the Boson and Fermion operators according to the time sequence assuming the commutation and anticommutation relations between these operators to be zero.

$$
\begin{align*}
T(\phi(\mathbf{x}) \phi(\mathbf{y})) & = \begin{cases}\phi(\mathbf{x}) \phi(\mathbf{y}), & \text { if } x_{0}>y_{0} \\
\phi(\mathbf{y}) \phi(\mathbf{x}), & \text { if } y_{0}>x_{0}\end{cases}  \tag{12.39}\\
T(\psi(\mathbf{x}) \psi(\mathbf{y})) & = \begin{cases}\psi(\mathbf{x}) \psi(\mathbf{y}), & \text { if } x_{0}>y_{0} \\
-\psi(\mathbf{y}) \psi(\mathbf{x}), & \text { if } y_{0}>x_{0}\end{cases}  \tag{12.40}\\
T(\psi(\mathbf{x}) \bar{\psi}(\mathbf{y})) & = \begin{cases}\psi(\mathbf{x}) \bar{\psi}(\mathbf{y}), & \text { if } x_{0}>y_{0} \\
-\bar{\psi}(\mathbf{y}) \psi(\mathbf{x}), & \text { if } y_{0}>x_{0}\end{cases} \tag{12.41}
\end{align*}
$$

In the above equations, the suffix 0 indicates that $x_{0}, y_{0}$ are the time components of the space-time coordinates $\mathbf{x}, \mathbf{y}$.

### 12.4.3 Wick's contractions

The basic problem is how to express Wick's chronological product in terms of a sum of normal products. For this, we introduce the notion of contraction between any two factors, known as Wick's contraction. The contraction represents the commutator or anticommutator bracket that arises from switching the factors in going from a chronological product to normal products.

$$
\begin{equation*}
T(A B)=A^{\bullet} B^{\bullet}+N(A B) \tag{12.42}
\end{equation*}
$$

where $A^{\bullet} B^{\bullet}$ represents the contracted pair. Some authors use the symbol of connecting bracket underneath the two operators to denote the contracted pair and write Eq. (12.42) as

$$
T(A B)=A B+N(A B),
$$

where $A B$ denotes the contracted pair, but we shall use the symbol $A^{\bullet} B^{\bullet}$ to denote Wick's contraction. Using Eq. (12.42) as the definition contracted pair, let us consider some simple cases and show that the contracted pair represents a commutation or anticommutation bracket which,
in turn, denotes a propagator.
For $x_{0}>y_{0}$,

$$
\begin{align*}
\phi^{+\bullet}(\mathbf{x}) \phi^{-\bullet}(\mathbf{y}) & =T\left(\phi^{+}(\mathbf{x}) \phi^{-}(\mathbf{y})\right)-N\left(\phi^{+}(\mathbf{x}) \phi^{-}(\mathbf{y})\right) \\
& =\phi^{+}(\mathbf{x}) \phi^{-}(\mathbf{y})-\phi^{-}(\mathbf{y}) \phi^{+}(\mathbf{x}) \\
& =\left[\phi^{+}(\mathbf{x}), \phi^{-}(\mathbf{y})\right]_{-} \\
& =i \hbar c \Delta^{+}(\mathbf{x}-\mathbf{y}) \tag{12.43}
\end{align*}
$$

For $y_{0}>x_{0}$,

$$
\begin{align*}
\phi^{+\bullet}(\mathbf{x}) \phi^{-\bullet}(\mathbf{y}) & =T\left(\phi^{+}(\mathbf{x}) \phi^{-}(\mathbf{y})\right)-N\left(\phi^{+}(\mathbf{x}) \phi^{-}(\mathbf{y})\right) \\
& =\phi^{-}(\mathbf{y}) \phi^{+}(\mathbf{x})-\phi^{-}(\mathbf{y}) \phi^{+}(\mathbf{x}) \\
& =0 \tag{12.44}
\end{align*}
$$

For $x_{0}>y_{0}$,

$$
\begin{align*}
\phi^{\bullet}(\mathbf{x}) \phi^{\bullet}(\mathbf{y}) & =T(\phi(\mathbf{x}) \phi(\mathbf{y}))-N(\phi(\mathbf{x}) \phi(\mathbf{y})) \\
& =\phi(\mathbf{x}) \phi(\mathbf{y})-N(\phi(\mathbf{x}) \phi(\mathbf{y})) \\
& =i \hbar c \Delta^{+}(\mathbf{x}-\mathbf{y}), \quad \text { using Eqn. (12.32). } \tag{12.45}
\end{align*}
$$

For $y_{0}>x_{0}$,

$$
\begin{align*}
\phi^{\bullet}(\mathbf{x}) \phi^{\bullet}(\mathbf{y}) & =T(\phi(\mathbf{x}) \phi(\mathbf{y}))-N(\phi(\mathbf{x}) \phi(\mathbf{y})) \\
& =\phi(\mathbf{y}) \phi(\mathbf{x})-N(\phi(\mathbf{x}) \phi(\mathbf{y})) \\
& =i \hbar c \Delta^{+}(\mathbf{y}-\mathbf{x}), \quad \text { using Eqn. (12.32). } \tag{12.46}
\end{align*}
$$

In a similar way, we can obtain the results for Fermion contracted pairs. For $x_{0}>y_{0}$,

$$
\begin{align*}
\psi_{\alpha}^{\bullet}(\mathbf{x}) \bar{\psi}_{\beta}^{\bullet}(\mathbf{y}) & =T\left(\psi_{\alpha}(\mathbf{x}) \bar{\psi}_{\beta}(\mathbf{y})\right)-N\left(\psi_{\alpha}(\mathbf{x}) \bar{\psi}_{\beta}(\mathbf{y})\right) \\
& =\psi_{\alpha}(\mathbf{x}) \bar{\psi}_{\beta}(\mathbf{y})-N\left(\psi_{\alpha}(\mathbf{x}) \bar{\psi}_{\beta}(\mathbf{y})\right) \\
& =-i S_{\alpha \beta}^{+}(\mathbf{x}-\mathbf{y}), \quad \text { using Eqn. (12.34) } \tag{12.47}
\end{align*}
$$

For $y_{0}>x_{0}$,

$$
\begin{align*}
\psi_{\alpha}^{\bullet}(\mathbf{x}) \bar{\psi}_{\beta}^{\bullet}(\mathbf{y}) & =T\left(\psi_{\alpha}(\mathbf{x}) \bar{\psi}_{\beta}(\mathbf{y})\right)-N\left(\psi_{\alpha}(\mathbf{x}) \bar{\psi}_{\beta}(\mathbf{y})\right) \\
& =-\bar{\psi}_{\beta}\left(\mathbf{y} \psi_{\alpha}(\mathbf{x})\right)-N\left(\psi_{\alpha}(\mathbf{x}) \bar{\psi}_{\beta}(\mathbf{y})\right) \\
& =i S_{\alpha \beta}^{-}(\mathbf{x}-\mathbf{y}), \quad \text { using Eqn. (12.33). } \tag{12.48}
\end{align*}
$$

Equations (12.47) and (12.48) together yield the Feynman propagator $K_{F}(\mathbf{x}-\mathbf{y})$.

$$
\begin{equation*}
\psi_{\alpha}^{\bullet}(\mathbf{x}) \bar{\psi}_{\beta}^{\bullet}(\mathbf{y})=K_{F_{\alpha \beta}}(\mathbf{x}-\mathbf{y}) . \tag{12.49}
\end{equation*}
$$

### 12.4.4 Wick's theorem

Wick has generalized the Wick's contraction formula (12.42) for a product of a large number of field operators $A\left(\mathbf{x}_{a}\right) B\left(\mathbf{x}_{b}\right) C\left(\mathbf{x}_{c}\right) \cdots Z\left(\mathbf{x}_{z}\right)$.

$$
\begin{align*}
& T(A B C D \cdots W X Y Z)=N(A B C D \cdots W X Y Z) \\
& \quad+N\left(A^{\bullet} B^{\bullet} C \cdots Z\right)+N\left(A^{\bullet} B C^{\bullet} \cdots Z\right)+\cdots+N\left(A \cdots X Y^{\bullet} Z^{\bullet}\right) \\
& \quad+N\left(A^{\bullet} B^{\bullet} C^{\bullet \bullet} D^{\bullet \bullet} \cdots W X Y Z\right)+\cdots+N\left(A B C D \cdots W^{\bullet} X^{\bullet} Y^{\bullet \bullet} Z^{\bullet \bullet}\right) \\
& \quad+\cdots+N\left(A^{\bullet} B^{\bullet} \cdots Y^{\bullet} \cdots Z^{\bullet} \bullet\right) \tag{12.50}
\end{align*}
$$

The right hand side of Eq. (12.50) consists of normal products with no contractions, with all possible contractions of one pair, two pairs, three pairs and so on. Different contractions are distinguished by using single dot, two dots, three dots and so on as superscripts. Equation (12.50) is known as Wick's theorem. It is proved by the method of induction: If it is true for a product of $n$ operators, it can be shown to be true for $n+1$ operators.

The contractions are mere numbers and not operators. In Eq. (12.50), although they are shown within the normal-ordered products, they can be pulled out of the normal ordered products.

The $S$-matrix expansion consists of time-ordered product of normalordered operators.

$$
\begin{equation*}
T\left(\mathscr{H}_{I}\left(\mathbf{x}_{1}\right) \cdots \mathscr{H}_{I}\left(\mathbf{x}_{n}\right)\right)=T\left(N(A B \cdots)_{\mathbf{x}_{1}} \cdots N\left((A B \cdots)_{\mathbf{x}_{n}}\right) .\right. \tag{12.51}
\end{equation*}
$$

Wick's theorem (12.50) can be extended to include such mixed T-products as shown in Eq. (12.51). In order to avoid commutation or anticommutation between the operators at the same space-time points, the creation operators are given a small increment in time $\epsilon$ over the annihilation operators so that the time ordering in each term coincides with normal ordering. At the end, $\epsilon$ can be put to zero. This yields

$$
\begin{align*}
& T\left(N(A B \cdots) \mathbf{x}_{1} \cdots N\left((A B \cdots) \mathbf{x}_{n}\right)\right. \\
& \quad=T\left((A B \cdots)_{\mathbf{x}_{1}} \cdots\left((A B \cdots)_{\mathbf{x}_{n}}\right)_{\text {noe.t.c. }}\right. \tag{12.52}
\end{align*}
$$

where the right hand side is expanded using Wick's theorem (12.50) with no equal-time contractions (no e.t.c.).

### 12.5 From S-matrix expansion to Feynman diagrams

The foregoing discussion on the reduction of the S-matrix into normal products using Wick's theorem will become clearer by considering specific examples in QED. Given the initial and final states, how to pick out the relevant terms in the S-matrix expansion of a given order to write down the transition amplitudes and draw the corresponding Feynman diagrams.

### 12.5.1 Interaction with external electromagnetic field

Given the interaction term

$$
\mathscr{H}_{I}(\mathbf{x})=j^{\mu}(\mathbf{x}) A_{\mu}^{e}(\mathbf{x})=\bar{\psi}(\mathbf{x}) \gamma^{\mu} \psi(\mathbf{x}) A_{\mu}^{e}(\mathbf{x})=\bar{\psi}(\mathbf{x}) A^{e}(\mathbf{x}) \psi(\mathbf{x})
$$

we can write down the S -matrix expansion as

$$
\begin{align*}
& S= 1+\left(\frac{-i e}{\hbar c}\right) \int d^{4} x_{1} N\left(\bar{\psi}(\mathbf{x}) A^{e}(\mathbf{x}) \psi(\mathbf{x})\right) \\
&+\left(\frac{-i e}{\hbar c}\right)^{2} \frac{1}{2!} \int d^{4} x_{1} d^{4} x_{2} T\left\{N\left(\bar{\psi}\left(\mathbf{x}_{1}\right) A^{e}\left(\mathbf{x}_{1}\right) \psi\left(\mathbf{x}_{1}\right)\right)\right. \\
&\left.\times N\left(\bar{\psi}\left(\mathbf{x}_{2}\right) A^{e}\left(\mathbf{x}_{2}\right) \psi\left(\mathbf{x}_{2}\right)\right)\right\}+\cdots . \tag{12.53}
\end{align*}
$$

## First order term in S-matrix

Let us consider the first order term in normal order form explicitly in terms of the positive and negative frequency (energy) parts $\psi^{-}(\mathbf{x})$ and $\psi^{+}(\mathbf{x})$.

$$
\begin{align*}
N\left(\bar{\psi}(\mathbf{x}) A^{e}(\mathbf{x}) \psi(\mathbf{x})\right)= & N\left\{\left(\bar{\psi}^{+}(\mathbf{x})+\bar{\psi}^{-}(\mathbf{x})\right) \mathcal{A}^{e}(\mathbf{x})\left(\psi^{+}(\mathbf{x})+\psi^{-}(\mathbf{x})\right\}\right. \\
= & \bar{\psi}^{+}(\mathbf{x}) \mathcal{A}^{e}(\mathbf{x}) \psi^{+}(\mathbf{x})-\psi^{-}(\mathbf{x}) \mathcal{A}^{e}(\mathbf{x}) \bar{\psi}^{+}(\mathbf{x}) \\
& +\bar{\psi}^{-}(\mathbf{x}) \mathcal{A}^{e}(\mathbf{x}) \psi^{+}(\mathbf{x})+\bar{\psi}^{-}(\mathbf{x}) \mathcal{A}^{e}(\mathbf{x}) \psi^{-}(\mathbf{x}) \cdot(12 \tag{12.54}
\end{align*}
$$

Let us use the Feynman convention of representing electron by an arrow pointing upwards (forward in time) and positron by an arrow pointing downwards (backward in time) as shown in Fig. 12.1. The positive and negative frequency parts of $\psi(\mathbf{x})$ and $\bar{\psi}(\mathbf{x})$ denote the following events:

$$
\begin{array}{lll}
\psi^{+}(\mathbf{x}) & : & \text { annihilation of electron } \\
\psi^{-}(\mathbf{x}) & : & \text { creation of positron } \\
\bar{\psi}^{+}(\mathbf{x}) & : & \text { annihilation of positron } \\
\bar{\psi}^{-}(\mathbf{x}) & : & \text { creation of electron }
\end{array}
$$



Figure 12.1: Pictorial representation of $\psi^{+}(\mathbf{x}), \psi^{-}(\mathbf{x}), \bar{\psi}^{+}(\mathbf{x}), \bar{\psi}^{-}(\mathbf{x})$.

Let us now represent the various terms in Eq. (12.54) by Feynman diagrams as shown in Fig. 12.2. They represent (a) the pair annihilation (b) the positron scattering (c) the electron scattering and (d) the pair creation. These are the basic vertices and all the real processes involving photons, electrons and positrons are built out of them. Energy and momentum are conserved at each vertex but the relativistic energymomentum relation for all the three particles participating in a vertex will not be simultaneously satisfied. So, all the three particles participating in a vertex will not be real and at least one of them has to be virtual for which the energy-momentum relation need not be satisfied. However, if the external electromagnetic interaction is the Coulomb potential (virtual photon), then the processes (b) and (c) representing positron scattering and electron scattering are possible. For any real process, there must be over-all energy-momentum conservation besides conservation of energy and momentum at each vertex.

We shall now consider the positron scattering by Coulomb potential as an example of the first-order process and show that we get the same scattering amplitude as obtained earlier by Feynman's method. The initial state $|i\rangle$ is specified by a positron of momentum $\boldsymbol{p}_{i}$ and spin state $s_{i}$ and the final state $|f\rangle$ by a positron of momentum $\boldsymbol{p}_{f}$ and spin state $s_{f}$.

$$
|i\rangle=d_{s_{i}}^{\dagger}\left(\boldsymbol{p}_{i}\right)|0\rangle ; \quad|f\rangle=d_{s_{f}}^{\dagger}\left(\boldsymbol{p}_{f}\right)|0\rangle
$$

The matrix element for the first order transition from an initial state to a final state is given by

$$
\begin{align*}
\mathcal{M} & =\langle f| S^{(1)}|i\rangle \\
& =-\frac{i e}{\hbar c} \int d^{4} x_{1}\langle f| N\left(\bar{\psi}\left(\mathbf{x}_{1}\right) \not A^{e}\left(\mathbf{x}_{1}\right) \psi\left(\mathbf{x}_{1}\right)\right)|i\rangle \\
& =-\frac{i e}{\hbar c} \int d^{4} x_{1}\langle 0| d_{s_{f}} N\left(\bar{\psi}\left(\mathbf{x}_{1}\right) A^{e}\left(\mathbf{x}_{1}\right) \psi\left(\mathbf{x}_{1}\right)\right) d_{s_{i}}^{\dagger}|0\rangle \tag{12.55}
\end{align*}
$$


$\bar{\psi}^{+}(\mathbf{x}) \AA^{e}(\mathbf{x}) \psi^{+}(\mathbf{x})$
(a) Pair annihilation

(b) Positron scattering


$$
\bar{\psi}^{-}(\boldsymbol{x}) \mathcal{A}^{e}(\mathbf{x}) \psi^{+}(\mathrm{x})
$$

(c) Electron scattering

(d) Pair creation

Figure 12.2: Feynman diagrams representing basic vertices in QED.

In the expansion of the normal product $N\left(\bar{\psi}\left(\mathbf{x}_{1}\right) \mathcal{A}^{e}\left(\mathbf{x}_{1}\right) \psi\left(\mathbf{x}_{1}\right)\right)$, only the term

$$
-\psi_{\beta}^{-}\left(\mathrm{x}_{1}\right) \mathcal{A}^{e} \bar{\psi}_{\alpha}^{+}\left(\mathrm{x}_{1}\right)
$$

will contribute to the process of positron scattering by an external potential $A^{e}$. The operator $\bar{\psi}_{\alpha}^{+}\left(\mathbf{x}_{1}\right)$ will destroy the positron in the initial state and $\psi_{\beta}^{-}\left(\mathbf{x}_{1}\right)$ will create the positron in the final state. Since

$$
\begin{align*}
\psi_{\beta}^{-}\left(\mathbf{x}_{1}\right) & =\frac{1}{(2 \pi)^{3 / 2}} \int_{p_{0}} d^{3} p \frac{1}{\sqrt{2 E_{p}}} \sum_{r=1}^{2} d_{r}^{\dagger}(\boldsymbol{p}) v_{r}(\boldsymbol{p}) e^{i \mathbf{p} \cdot \mathbf{x}_{1}},  \tag{12.56}\\
\bar{\psi}_{\alpha}^{+}\left(\mathbf{x}_{1}\right) & =\frac{1}{(2 \pi)^{3 / 2}} \int_{p_{0}^{\prime}} d^{3} p^{\prime} \frac{1}{\sqrt{2 E_{p^{\prime}}}} \sum_{s=1}^{2} d_{s}\left(\boldsymbol{p}^{\prime}\right) \bar{v}_{s}\left(\boldsymbol{p}^{\prime}\right) e^{-i \mathbf{p}^{\prime} \cdot \mathbf{x}_{1}}, \tag{12.57}
\end{align*}
$$



Figure 12.3: Scattering of positron by an external Coulomb field $\phi . \boldsymbol{p}_{i}$ and $\boldsymbol{p}_{f}$ denote the momenta of the incoming and scattered positron and $\boldsymbol{q}$, the threemomentum transfer to the positron.
we get

$$
\begin{align*}
\mathcal{M}= & \frac{i e}{\hbar c} \frac{1}{(2 \pi)^{3}} \int d^{4} x_{1} \int d^{3} p \int d^{3} p^{\prime} \sum_{r, s}^{1,2}\langle 0| d_{s_{f}}\left(\boldsymbol{p}_{f}\right) d_{r}^{\dagger}(\boldsymbol{p}) d_{s}\left(\boldsymbol{p}^{\prime}\right) d_{s_{i}}^{\dagger}\left(\boldsymbol{p}_{i}\right)|0\rangle \\
& \times \frac{1}{\sqrt{2 E_{p} 2 E_{p^{\prime}}}} \bar{v}_{s}\left(\boldsymbol{p}^{\prime}\right) \mathcal{A}^{e}\left(\mathbf{x}_{1}\right) v_{r}(\boldsymbol{p}) e^{i\left(\mathbf{p}-\mathbf{p}^{\prime}\right) \cdot \mathbf{x}_{1}} . \tag{12.58}
\end{align*}
$$

It can be easily verified that

$$
\begin{equation*}
\langle 0| d_{s_{f}}\left(\boldsymbol{p}_{f}\right) d_{r}^{\dagger}(\boldsymbol{p}) d_{s}\left(\boldsymbol{p}^{\prime}\right) d_{s_{i}}^{\dagger}\left(\boldsymbol{p}_{i}\right)|0\rangle=\delta_{r s_{f}} \delta_{s s_{i}} \delta\left(\boldsymbol{p}-\boldsymbol{p}_{f}\right) \delta\left(\boldsymbol{p}^{\prime}-\boldsymbol{p}_{i}\right) \tag{12.59}
\end{equation*}
$$

Substituting (12.59) into Eq. (12.58) and integrating over $d^{3} p$ and $d^{3} p^{\prime}$, we obtain

$$
\begin{equation*}
\mathcal{M}=\frac{i e}{\hbar c} \frac{1}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p_{i}} 2 E_{p_{f}}}} \int d^{4} x_{1} \bar{v}_{s_{i}}\left(\boldsymbol{p}_{i}\right) \not A^{e}\left(\mathbf{x}_{1}\right) v_{s_{f}}\left(\boldsymbol{p}_{f}\right) e^{i\left(\mathbf{p}_{f}-\mathbf{p}_{i}\right) \cdot \mathbf{x}_{1}} .( \tag{12.60}
\end{equation*}
$$

Using the Fourier transform relation

$$
d(\mathbf{q})=\frac{1}{(2 \pi)^{4}} \int A^{e}\left(\mathbf{x}_{1}\right) e^{i \mathbf{q} \cdot \mathbf{x}} d^{4} x
$$

we obtain the matrix element in the momentum representation.

$$
\begin{equation*}
\mathcal{M}=\left(\frac{i e}{\hbar c}\right) \frac{2 \pi}{\sqrt{2 E_{p_{i}} 2 E_{p_{f}}}} \bar{v}_{s_{i}}\left(\boldsymbol{p}_{i}\right) d\left(\mathbf{p}_{f}-\mathbf{p}_{i}\right) v_{s_{f}}\left(\boldsymbol{p}_{f}\right) . \tag{12.61}
\end{equation*}
$$

This is exactly the matrix element that one would have obtained in Feynman's theory. From the equation obtained for the matrix element,
one can frame the rule for drawing the Feynman diagram in momentum space for positron scattering. The initial and final states follow the path sequence. If we are considering positron scattering, the final state occurs at the bottom of the diagram (with time flowing upwards) as indicated in Fig. 12.3.

Thus we have established the connection between the field theoretical formalism and the much simpler Feynman formalism developed intuitively by him.

### 12.5.2 Second order term in S-matrix

Let us now consider the second order term in the S-matrix and show how helpful is Wick's theorem in expanding it into a sum of normal products which can be used to study a large variety of quantum electrodynamic processes. The second order term involves a product of two $\mathscr{H}_{I}(\mathbf{x})$ and the basic vertex diagrams which we obtained in the first order can be combined in a meaningful way to yield Feynman diagrams that represent real processes.

Using Wick's theorem, let us expand the second order term in the S-matrix as a sum of normal products.

$$
\begin{align*}
T & \left(\mathscr{H}_{I}\left(\mathbf{x}_{1}\right) \mathscr{H}_{I}\left(\mathbf{x}_{2}\right)\right) \\
& =T\left\{N\left(\bar{\psi}\left(\mathbf{x}_{1}\right) A\left(\mathbf{x}_{1}\right) \psi\left(\mathbf{x}_{1}\right)\right) N\left(\bar{\psi}\left(\mathbf{x}_{2}\right) A\left(\mathbf{x}_{2}\right) \psi\left(\mathbf{x}_{2}\right)\right)\right\} \\
& =N_{A}+N_{B}+N_{C}+N_{D}+N_{E}+N_{F}+N_{G}+N_{H} \tag{12.62}
\end{align*}
$$

where

$$
\begin{align*}
& N_{A}=N\left(\bar{\psi}\left(\mathbf{x}_{1}\right) A\left(\mathbf{x}_{1}\right) \psi\left(\mathbf{x}_{1}\right) \bar{\psi}\left(\mathbf{x}_{2}\right) A\left(\mathbf{x}_{2}\right) \psi\left(\mathbf{x}_{2}\right)\right),  \tag{12.63}\\
& N_{B}=N\left(\bar{\psi}\left(\mathbf{x}_{1}\right) A\left(\mathbf{x}_{1}\right) \psi \cdot\left(\mathbf{x}_{1}\right) \bar{\psi} \cdot\left(\mathbf{x}_{2}\right) A\left(\mathbf{x}_{2}\right) \psi\left(\mathbf{x}_{2}\right)\right),  \tag{12.64}\\
& N_{C}=N\left(\bar{\psi}\left(\mathbf{x}_{1}\right) A\left(\mathbf{x}_{1}\right) \psi\left(\mathbf{x}_{1}\right) \bar{\psi}\left(\mathbf{x}_{2}\right) A\left(\mathbf{x}_{2}\right) \psi \cdot\left(\mathbf{x}_{2}\right)\right),  \tag{12.65}\\
& N_{D}=N\left(\bar{\psi}\left(\mathbf{x}_{1}\right) \gamma^{\alpha} A_{\alpha}^{\bullet}\left(\mathbf{x}_{1}\right) \psi\left(\mathbf{x}_{1}\right) \bar{\psi}\left(\mathbf{x}_{2}\right) \gamma^{\alpha} A_{\alpha}^{\bullet}\left(\mathbf{x}_{2}\right) \psi\left(\mathbf{x}_{2}\right)\right),  \tag{12.66}\\
& N_{E}=N\left(\bar{\psi}\left(\mathbf{x}_{1}\right) \gamma^{\alpha} A_{\alpha}^{\bullet}\left(\mathbf{x}_{1}\right) \psi^{\bullet \bullet}\left(\mathbf{x}_{1}\right) \bar{\psi} \bullet \bullet\left(\mathbf{x}_{2}\right) \gamma^{\alpha} A_{\alpha}^{\bullet}\left(\mathbf{x}_{2}\right) \psi\left(\mathbf{x}_{2}\right)\right),  \tag{12.67}\\
& N_{F}=N\left(\bar{\psi}^{\bullet}\left(\mathbf{x}_{1}\right) \gamma^{\alpha} A_{\alpha}^{\bullet \bullet}\left(\mathbf{x}_{1}\right) \psi\left(\mathbf{x}_{1}\right) \bar{\psi}\left(\mathbf{x}_{2}\right) \gamma^{\alpha} A_{\alpha}^{\bullet \bullet}\left(\mathbf{x}_{2}\right) \psi^{\bullet}\left(\mathbf{x}_{2}\right)\right),  \tag{12.68}\\
& N_{G}=N\left(\bar{\psi} \bullet\left(\mathbf{x}_{1}\right) A\left(\mathbf{x}_{1}\right) \psi^{\bullet \bullet}\left(\mathbf{x}_{1}\right) \bar{\psi} \bullet \bullet\left(\mathbf{x}_{2}\right) A\left(\mathbf{x}_{2}\right) \psi^{\bullet \bullet}\left(\mathbf{x}_{2}\right)\right),  \tag{12.69}\\
& N_{H}=N\left(\bar{\psi}\left(\mathbf{x}_{1}\right) \gamma^{\alpha} A_{\alpha}^{\bullet \bullet}\left(\mathbf{x}_{1}\right) \psi^{\bullet \bullet \bullet}\left(\mathbf{x}_{1}\right) \bar{\psi} \bullet \bullet \bullet\left(\mathbf{x}_{2}\right) \gamma^{\alpha} A_{\alpha}^{\bullet \bullet}\left(\mathbf{x}_{2}\right) \psi^{\bullet}\left(\mathbf{x}_{2}\right)\right) . \tag{12.70}
\end{align*}
$$

The first term $N_{A}$ is not very interesting. It corresponds to two disconnected vertex diagrams and does not represent any real process.

The second and third terms $N_{B}$ and $N_{C}$ give identical results since permutation of Fermion operators in the two groups $(\bar{\psi} A \psi)_{\mathbf{x}_{1}}$ and $(\bar{\psi} A \psi)_{\mathbf{x}_{2}}$
is even and the space-time variables are dummy variables that are integrated. Expressions (12.64) and (12.65) contain one Fermion contraction which represents the Fermion propagator from one space-time point to another besides two uncontracted Fermion operators and two uncontracted photon operators. The uncontracted operators annihilate the particles that are present initially and emit the particles that are present finally. Since the operators are arranged in normal order, it is easy to pick out the terms that contribute to a given process. The terms $N_{B}$ and $N_{C}$ can describe several processes involving two Fermions and two photons such as Compton scattering by electron, Compton scattering by positron, electron-positron pair annihilation into two photons and electron-positron pair creation by two photons.

$$
\begin{array}{lllll}
\text { (i) } \gamma+e^{-} & \rightarrow \gamma+e^{-}, & \text {(ii) } \gamma+e^{+} & \rightarrow \gamma+e^{+}, \\
\text {(iii) } e^{+}+e^{-} & \rightarrow \gamma+\gamma, & \text { (iv) } \gamma+\gamma & \rightarrow e^{+}+e^{-} .
\end{array}
$$

The fourth term $N_{D}$ consists of four uncontracted Fermion operators and one photon-photon contraction representing the photon propagation between the two Fermions. By picking out suitable Fermion annihilation and creation operators from the uncontracted Fermion operators, we can describe electron-electron scattering, electron-positron scattering and positron-positron scattering. For each of these processes, there will be two Feynman diagrams. For electron-electron and positron-positron scattering, there will be two Feynman amplitudes, one for direct scattering and another for exchange scattering with a relative change in sign. In non-relativistic quantum mechanics, the relative minus sign arises from the use of antisymmetric wave function of Fermions. In the field theory, this is provided by the anticommutative property of the Fermion field operators. In the case of electron-positron scattering (known as Bhabha scattering), there are also two Feynman diagrams, one direct scattering and the other representing electron-positron pair annihilation into a photon and subsequent pair creation. These two modes are indistinguishable and should be taken into account in the calculation.

The fifth and sixth terms $N_{E}$ and $N_{F}$ can be clubbed together. Each term consists of contraction of one Fermion pair and one photon-pair, yielding one Fermion propagator and one photon propagator. The electron emits a photon which is recaptured by the same electron. This is known as the self-energy of the electron and is represented by the Feynman diagram 12.4(a). This process converts the bare electron into a physical electron. The self-interaction energy changes the mass of the bare elec-
tron into that of the physical electron. Its evaluation leads to a divergent integral which is overcome by the method of renormalization.


Figure 12.4: Feynman diagrams representing (a) Self energy of the electron (b) Self energy of the photon or Vacuum polarization and (c) Vacuum diagram.

The seventh term $N_{G}$ consists of two Fermion pair contractions and two uncontracted photon field operators. The photon creates a virtual electron-positron pair which again annihilates into a photon. This is known as the self-energy of the photon and is depicted by the Feynman diagram (b) in Fig. 12.4. This is also known as vacuum polarization since the external electromagnetic field such as that of heavy nucleus can modify the distribution of virtual electron-positron pairs in the vacuum. The calculation of the photon self-energy also leads to infinities which are again eliminated by renormalization procedure.

The last term $N_{H}$ is one in which the two Fermion-pairs and one photon-pair are contracted leaving no external line. This is a vacuum diagram as illustrated in Fig. 12.4(c). Such vacuum diagrams without any external line can be omitted.

## Review Questions

12.1 Develop an invariant perturbation theory treating the interaction between fields as a small perturbation and deduce the S-matrix formalism for the study of scattering phenomena.
12.2 Define Normal Products, Dyson's Chronological products, Wick's chronological products and Wick's contraction. State Wick's theorem and explain how it is useful in reducing Wick's chronological product into normal products.
12.3 Consider the first order term in the S-matrix expansion for the interaction of the Dirac field with an external electromagnetic field and explain how you can represent the various processes by means of Feynman diagrams. Deduce the Feynman matrix element in momentum representation for the positron scattering by an external electromagnetic field.
12.4 Explain how the S-matrix can be decomposed into normal products using Wick's theorem. Illustrate your answer by considering the second order term in the S-matrix for the electro-magnetic interaction.

## Problems

12.1 Given Eq. (12.3)

$$
\left(H_{0}+H_{i}\right) \Phi(t)=i \hbar \frac{\partial \Phi(t)}{\partial t}
$$

where $H_{0}$ and $H_{i}$ are the unperturbed and interaction Hamiltonian and $\Phi(t)$ is the time-dependent state vector in the Schrödinger picture, deduce Eq. (12.5)

$$
i \hbar \frac{\partial \Psi(t)}{\partial t}=H_{I}(t) \Psi(t)
$$

in the interaction picture where the state vector $\Psi(t)$ and the interaction Hamiltonian $H_{I}(t)$ are defined by

$$
\Psi(t)=e^{i H_{0} t / \hbar} \Phi(t) \quad \text { and } \quad H_{I}(t)=e^{i H_{0} t / \hbar} H_{i} e^{-i H_{0} t / \hbar}
$$

10.2 Given the Dirac spinors and the gamma matrices, construct a scalar, a pseudoscalar, a four-vector, a pseudo vector, and an antisymmetric tensor of second rank. With the help of these, construct an interaction Hamiltonian for electron-photon interaction and nucleon-pion interaction.
12.3 It is known that parity is not conserved in weak interaction. So, construct an interaction Hamiltonian for the parity-violating $\beta$-decay.

## Solutions to Problems

12.1 Since $\Psi(t)=e^{i H_{0} t / \hbar} \Phi(t)$ and $i \hbar \frac{\partial \Phi(t)}{\partial t}=\left(H_{0}+H_{i}\right) \Phi(t)$, we get by differentiating $\Psi$,

$$
\begin{aligned}
\frac{\partial \Psi}{\partial t} & =\frac{i}{\hbar} H_{0} e^{i H_{0} t / \hbar} \Phi(t)+e^{i H_{0} t / \hbar} \frac{\partial \Phi}{\partial t} \\
& =\frac{i}{\hbar} H_{0} e^{i H_{0} t / \hbar} \Phi(t)+e^{i H_{0} t / \hbar} \frac{1}{i \hbar}\left(H_{0}+H_{i}\right) \Phi(t) \\
& =\frac{1}{i \hbar} e^{i H_{0} t / \hbar} H_{i} \Phi(t) \\
& =\frac{1}{i \hbar} e^{i H_{0} t / \hbar} H_{i} e^{-i H_{0} t / \hbar} \Psi(t) .
\end{aligned}
$$

Hence it follows that

$$
i \hbar \frac{\partial \Psi}{\partial t}=H_{I}(t) \Psi(t),
$$

where $H_{I}(t)=e^{i H_{0} t / \hbar} H_{i} e^{-i H_{0} t / \hbar}$.
12.2 Given the Dirac spinors and $\gamma$ matrices, we can construct the following quantities which have definite transformation property:

| Quantity | Transformation Property | No. of components |
| :---: | :---: | :---: |
| $\bar{\psi} \psi$ | Scalar | 1 |
| $\bar{\psi} \gamma_{5} \psi$ | Pseudoscalar | 1 |
| $i \bar{\psi} \gamma^{\mu} \psi$ | Four-vector | 4 |
| $i \bar{\psi} \gamma_{\mu} \gamma_{5} \psi$ | Pseudo-vector | 4 |
| $\frac{i}{2} \bar{\psi}\left(\gamma^{\mu} \gamma^{\nu}-\gamma^{\nu} \gamma^{\mu}\right) \psi$ | Antisymmetric tensor | 6 |
|  | (2 ${ }^{\text {nd }}$ rank $)$ |  |

The factor $i$ is included in order to take care of the Hermitian property.
Since the photon field is represented by a four-vector $A_{\mu}$, the electronphoton interaction is represented by the interaction Hamiltonian

$$
H_{\mathrm{int}}=\bar{\psi} \gamma^{\mu} \psi A_{\mu},
$$

which is a scalar under Lorentz transformation.
The pion-nucleon interaction is a bit complicated. Since pion has an intrinsic negative parity when compared to the nucleon, the interaction Hamiltonian will have a form

$$
\bar{\psi}_{N} \gamma_{5} \psi_{N} \phi
$$

where $\phi$ represents the pseudo-scalar pion field. This interaction Hamiltonian is a scalar under both Lorentz transformation and parity operation. Since the pion exists in three charge states, it can be represented by an isospin vector $\boldsymbol{\tau}_{\pi}$. Since the nucleon exists in two charge states, it can be represented by a vector $\boldsymbol{\tau}_{N}$ in isospin space. The interaction Hamiltonian has to be a scalar in the iso-spin space too. This yield an interaction Hamiltonian

$$
\bar{\psi}_{N} \gamma_{5} \boldsymbol{\tau}_{N} \cdot \boldsymbol{\tau}_{\pi} \psi_{N} \phi_{\pi}
$$

where $\psi_{N}$ is a two-component vector and $\phi_{\pi}$ is a three-component vector in the isospin space.

$$
\psi_{N}=\left[\begin{array}{c}
\psi_{p} \\
\psi_{n}
\end{array}\right] ; \quad \phi_{\pi}=\left[\begin{array}{c}
\phi_{1} \\
\phi_{2} \\
\phi_{3}
\end{array}\right] .
$$

12.3 It is known that the beta decay interaction is a combination of vector and axial vector interaction.

$$
H_{\mathrm{int}}=G_{\beta} \sum_{i=V, A} C_{i}\left(\bar{\psi}_{p} \mathcal{O}_{i} \psi_{n}\right)\left(\bar{\psi}_{e} \mathcal{O}_{i} \psi_{\nu}\right),
$$

where $G_{\beta}$ is the weak interaction coupling constant and $C_{V}$ and $C_{A}$ are the coupling coefficients for vector and axial vector interaction. The operators $\mathcal{O}_{V}=\gamma^{\mu}$ and $\mathcal{O}_{A}=\gamma^{\mu} \gamma_{5}$ This is a parity conserving interaction. Subsequently Lee and Yang postulated that weak interaction need not be conserved in weak interaction and it has been experimentally confirmed by Wu et al. It is found that $C_{A}=-C_{V}$ for pure leptonic weak interaction such as $\mu$-decay. Incorporating this in $V-A$ interaction, we get

$$
H_{\mathrm{int}}=\frac{G_{\beta}}{\sqrt{2}} C_{V}\left[\bar{\psi}_{p} \gamma^{\mu}\left(1-\gamma_{5}\right) \psi_{n}\right]\left[\bar{\psi}_{e} \gamma_{\mu}\left(1-\gamma_{5}\right) \psi_{\nu}\right]
$$

The inclusion of the parity non-conserving extra term requires the insertion of a factor $1 / \sqrt{2}$ in order to retain the old definition of the coupling constant $G_{\beta}$.

## Chapter 13

## The Gauge Theories

The Lagrangian density $\mathscr{L}$ of a field should be invariant under a phase transformation of the field variable. This is often referred to as the invariance of the Lagrangian density under gauge transformation. If the phase transformation is independent of the space-time coordinates, then it is known as the global phase transformation. The invariance of $\mathscr{L}$ under global phase transformation implies certain symmetry property and leads to the conservation of current. This in turn, leads to the law of conservation of charge or conservation of lepton number or baryon number.

If the phase transformation depends on the space-time coordinates $(\mathbf{x})$, then it is called the local phase transformation or the local gauge transformation. In order to make $\mathscr{L}$ invariant under local gauge transformation, certain additional vector fields known as gauge fields have to be introduced. The gauge fields specify the interaction dynamics among the fields. Thus gauge theories are of special interest since they govern the interaction dynamics. The $U(1)$ gauge invariance leads to the electromagnetic interaction, the $S U(2) \times U(1)$ gauge invariance leads to the electro-weak interaction and $S U(3)$ gauge invariance leads to the strong interaction. The $U(1), S U(2) \times U(1)$ and $S U(3)$ are the symmetry groups that specify the gauge transformation.

One may ask why the Lagrangian density $\mathscr{L}$ should be invariant under gauge transformation. It is an important criterion that is to be fulfilled for developing a renormalizable field theory. That is why gauge theories are of paramount importance in the formulation of quantum field theory.

### 13.1 Invariance of Lagrangian under gauge transformation

## Noether's theorem

Consider the Dirac Lagrangian density of a fermion field

$$
\begin{equation*}
\mathscr{L}=\bar{\psi}(\mathbf{x})\left(i \gamma^{\mu} \partial_{\mu}-m\right) \psi(\mathbf{x}) . \tag{13.1}
\end{equation*}
$$

In quantum field theory, the Lagrangian density should be Hermitian. So, let us use the equivalent Dirac Lagrangian density which is Hermitian.

$$
\begin{equation*}
\mathscr{L}=\frac{i}{2}\left(\bar{\psi} \gamma^{\mu} \partial_{\mu} \psi-\left(\partial_{\mu} \bar{\psi}\right) \gamma^{\mu} \psi\right)-m \bar{\psi} \psi . \tag{13.2}
\end{equation*}
$$

It can be easily checked that the Lagrangian $\mathscr{L}$ is invariant under the phase transformation

$$
\begin{equation*}
\psi(\mathbf{x}) \rightarrow e^{i \alpha} \psi(\mathbf{x}) \tag{13.3}
\end{equation*}
$$

where $\alpha$ is a real constant. It follows that

$$
\begin{align*}
\partial_{\mu} \psi & \rightarrow e^{i \alpha} \partial_{\mu} \psi .  \tag{13.4}\\
\bar{\psi} & \rightarrow e^{-i \alpha} \bar{\psi} . \tag{13.5}
\end{align*}
$$

The phase transformation $U(\alpha)=e^{i \alpha}$, with a single parameter $\alpha$ running over all real numbers, forms a unitary Abelian group $U(1)$. It is called Abelian since the group multiplication is commutative.

$$
U\left(\alpha_{1}\right) U\left(\alpha_{2}\right)=U\left(\alpha_{2}\right) U\left(\alpha_{1}\right) .
$$

The invariance of $\mathscr{L}$ under unitary transformation $U(1)$ implies a symmetry property and according to Noether's theorem, any symmetry property implies a conservation law. To observe this, let us consider the invariance of $\mathscr{L}$ under an infinitesimal $U(1)$ transformation.

$$
\begin{equation*}
\psi \rightarrow(1+i \alpha) \psi ; \quad \bar{\psi} \rightarrow(1-i \alpha) \bar{\psi} \tag{13.6}
\end{equation*}
$$

Accordingly, we have

$$
\begin{align*}
\delta \psi=i \alpha \psi ; & \delta\left(\partial_{\mu} \psi\right)=i \alpha\left(\partial_{\mu} \psi\right) .  \tag{13.7}\\
\delta \bar{\psi}=-i \alpha \bar{\psi} ; & \delta\left(\partial_{\mu} \bar{\psi}\right)=-i \alpha\left(\partial_{\mu} \bar{\psi}\right) . \tag{13.8}
\end{align*}
$$

Invariance of $\mathscr{L}$ implies that

$$
\delta \mathscr{L}=\frac{\partial \mathscr{L}}{\partial \psi} \delta \psi+\frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \psi\right)} \delta\left(\partial_{\mu} \psi\right)+\delta \bar{\psi} \frac{\partial \mathscr{L}}{\partial \bar{\psi}}+\delta\left(\partial_{\mu} \bar{\psi}\right) \frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \bar{\psi}\right)}=0 .(
$$

In the Lagrangian (13.2), the mass term will not contribute anything to $\delta \mathscr{L}$. So, we need to consider only the first two terms in (13.2). The terms (1) and (2) in Eq. (13.9) yield

$$
\begin{align*}
(1)+(2) & =\frac{\partial \mathscr{L}}{\partial \psi}(i \alpha \psi)+\frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \psi\right)}\left(i \alpha \partial_{\mu} \psi\right) \\
& =i \alpha\left\{\frac{\partial \mathscr{L}}{\partial \psi}-\partial_{\mu}\left(\frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \psi\right)}\right)\right\} \psi+i \alpha \partial_{\mu}\left(\frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \psi\right)} \psi\right) \\
& =i \alpha \partial_{\mu}\left(\frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \psi\right)} \psi\right), \tag{13.10}
\end{align*}
$$

since the quantity within the curly bracket is zero due to the EulerLagrange relation.

In a similar way, we can evaluate the terms (3) and (4) in Eq. (13.9)

$$
\begin{align*}
(3)+(4) & =(-i \alpha \bar{\psi}) \frac{\partial \mathscr{L}}{\partial \bar{\psi}}-i \alpha\left(\partial_{\mu} \bar{\psi}\right) \frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \bar{\psi}\right)} \\
& =-i \alpha \bar{\psi}\left\{\frac{\partial \mathscr{L}}{\partial \bar{\psi}}-\partial_{\mu}\left(\frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \bar{\psi}\right)}\right)\right\}-i \alpha \partial_{\mu}\left(\bar{\psi} \frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \bar{\psi}\right)}\right) \\
& =-i \alpha \partial_{\mu}\left(\bar{\psi} \frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \bar{\psi}\right)}\right), \tag{13.11}
\end{align*}
$$

since the quantity within the curly bracket is zero due to the EulerLagrange relation.

Substituting (13.10) and (13.11) into Eq. (13.9), we get

$$
\begin{align*}
\delta \mathscr{L} & =i \alpha \partial_{\mu}\left[\frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \psi\right)} \psi-\bar{\psi} \frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \bar{\psi}\right)}\right] \\
& =\partial_{\mu}\left(-\alpha \bar{\psi} \gamma^{\mu} \psi\right)=0 . \tag{13.12}
\end{align*}
$$

The Lagrangian (13.2) is used to obtain Eq. (13.12). It represents the conserved current $\partial_{\mu} j^{\mu}=0$ with

$$
\begin{equation*}
j^{\mu}=-\alpha \bar{\psi} \gamma^{\mu} \psi . \tag{13.13}
\end{equation*}
$$

Thus the invariance of the Dirac Lagrangian under the global phase transformation leads to the conservation of current $j^{\mu}$ as per Noether's theorem.

Let us consider the zeroth component $j^{0}$ of the four-vector current $j^{\mu}$. It yields the law of conservation of charge.

$$
\begin{align*}
Q & =\int d^{3} x j^{0}=-\alpha \int d^{3} x \bar{\psi} \gamma^{0} \psi=-\alpha \int d^{3} x \psi^{\dagger} \psi \\
& =-\alpha \int d^{3} p \sum_{r}\left(c_{r}^{\dagger}(\boldsymbol{p}) c_{r}(\boldsymbol{p})-d_{r}^{\dagger}(\boldsymbol{p}) d_{r}(\boldsymbol{p})\right) \tag{13.14}
\end{align*}
$$

where $c_{r}, c_{r}^{\dagger}, d_{r}$ and $d_{r}^{\dagger}$ are the annihilation and creation operators of the Fermions and anti-Fermions. If $\alpha$ is the charge $e$, then it represents the conservation of charge. Eqs. (13.13) and (13.14) are exactly the same as Eqs. (10.32) and (10.33). The concept of conserved charge can be extended to include the conservation of lepton number or baryon number, since Eq. (13.14) refers to any property that depends on the number of Fermions minus the number of anti-Fermions.

### 13.2 Gauge theory of electromagnetic interaction

Let us investigate what happens if we impose the global phase rotation symmetry to be also a local symmetry. In this case, $\alpha$ is a function of space-time coordinates $x^{\mu}$. From Eq. (13.3), it follows

$$
\begin{align*}
\partial_{\mu} \psi^{\prime} & =\partial_{\mu}\left(e^{i \alpha(\mathbf{x})} \psi(\mathbf{x})\right) \\
& =\left(\partial_{\mu} e^{i \alpha}\right) \psi+e^{i \alpha} \partial_{\mu} \psi \\
& =i\left(\partial_{\mu} \alpha\right) e^{i \alpha} \psi+e^{i \alpha} \partial_{\mu} \psi . \tag{13.15}
\end{align*}
$$

For Fermions,

$$
\begin{align*}
\mathscr{L}_{\psi^{\prime}} & =\bar{\psi}^{\prime}(i \not \partial-m) \psi^{\prime} \\
& =-\bar{\psi}\left(\partial_{\mu} \alpha\right) \gamma^{\mu} \psi+\bar{\psi}(i \not \partial-m) \psi \\
& =-\left(\partial_{\mu} \alpha\right) \bar{\psi} \gamma^{\mu} \psi+\mathscr{L}_{\psi} . \tag{13.16}
\end{align*}
$$

Equation (13.16) clearly indicates that the Dirac Lagrangian is not invariant under local phase rotation. This failure can be rectified by introducing new fields such that the modified Lagrangian has the local symmetry. The extra term on the right-hand side of Eq. (13.16) involves a factor $\partial_{\mu} \alpha$
which transforms like a four-vector. Introducing an additional term with a four-vector field $A_{\mu}$ in the Dirac Lagrangian, a modified Lagrangian

$$
\begin{equation*}
\mathscr{L}=\bar{\psi}(i \not \partial-m) \psi+e \bar{\psi} \gamma^{\mu} \psi A_{\mu} \tag{13.17}
\end{equation*}
$$

can be obtained which has a local symmetry under the transformation

$$
\begin{equation*}
\psi \rightarrow \psi^{\prime}=e^{i \alpha(\mathbf{x})} \psi, \quad A_{\mu} \rightarrow A_{\mu}^{\prime} \tag{13.18}
\end{equation*}
$$

What should be the transformation property of $A_{\mu}$ such that the modified Lagrangian is invariant under local phase transformation? The modified Lagrangian $\mathscr{L}^{\prime}$

$$
\begin{align*}
\mathscr{L}^{\prime} & =\bar{\psi}^{\prime}(i \not \partial-m) \psi^{\prime}+e \bar{\psi}^{\prime} \gamma^{\mu} \psi^{\prime} A_{\mu}^{\prime} \\
& =\bar{\psi}(i \not \partial-m) \psi+e\left(A_{\mu}^{\prime}-\frac{1}{e} \partial_{\mu} \alpha\right) \bar{\psi} \gamma^{\mu} \psi \tag{13.19}
\end{align*}
$$

is invariant if the newly introduced four-vector field $A_{\mu}$ transforms as

$$
\begin{equation*}
A_{\mu}^{\prime}=A_{\mu}+\frac{1}{e} \partial_{\mu} \alpha . \tag{13.20}
\end{equation*}
$$

Now one can easily identify that the newly introduced field $A_{\mu}$ is none other than the electromagnetic field which leaves the Lagrangian of the free electromagnetic field invariant under the same local transformation. Adding the free field Lagrangian of the electromagnetic field ${ }^{1}$, we obtain the full Lagrangian for the combined field of $\psi$ and $A_{\mu}$.

$$
\begin{equation*}
\mathscr{L}=\bar{\psi}(i \not \partial-m) \psi-\frac{1}{4} F_{\mu \nu} F^{\mu \nu}+e \bar{\psi} \gamma^{\mu} \psi A_{\mu} \tag{13.21}
\end{equation*}
$$

This is the Lagrangian that is used in field theory for the study of Quantum Electrodynamics (QED). It may be observed that the full Lagrangian (13.21) consists of an interaction term

$$
\mathscr{L}_{\mathrm{int}}=e \bar{\psi} \gamma^{\mu} \psi A_{\mu}
$$

besides the free-field Lagrangians of Dirac and e.m. fields. Since the origin of the interaction term can be traced to the introduction of a new field to make the Dirac Lagrangian invariant under a local gauge transformation, the newly introduced field has come to be known as gauge field, its field quantum as gauge Boson and the underlying theory as gauge theory.

[^74]is gauge invariant, the total Lagrangian $\mathscr{L}$ remains gauge invariant.

## Gauge covariant derivative

Let us define a gauge covariant derivative ${ }^{2} D_{\mu}$.

$$
\begin{equation*}
D_{\mu}=\partial_{\mu}-i e A_{\mu} \tag{13.22}
\end{equation*}
$$

such that the complete Lagrangian for QED (Eq. (13.21)) can be written in the form

$$
\begin{equation*}
\mathscr{L}=\bar{\psi}\left(i \gamma^{\mu} D_{\mu}-m\right) \psi-\frac{1}{4} F_{\mu \nu} F^{\mu \nu} . \tag{13.23}
\end{equation*}
$$

The same technique can be used for any other field with global symmetry. Let us illustrate it for the charged scalar field $\phi$. Any complex scalar field $\phi$ which has a global phase symmetry for the transformation

$$
\phi \rightarrow e^{i \alpha} \phi
$$

will also acquire a local phase symmetry with the substitution

$$
\partial_{\mu} \rightarrow D_{\mu} ; \quad A_{\mu} \rightarrow A_{\mu}+\frac{1}{e} \partial_{\mu} \alpha
$$

Thus we can obtain the Lagrangian for the scalar field interacting with the electromagnetic field as

$$
\begin{equation*}
\mathscr{L}=\left(D_{\mu} \phi\right)^{\dagger}\left(D_{\mu} \phi\right)-m^{2} \phi^{\dagger} \phi-\frac{1}{4} F_{\mu \nu} F^{\mu \nu} . \tag{13.24}
\end{equation*}
$$

### 13.3 Spontaneous symmetry breaking

The Lagrangian density of any field has a certain symmetry. Under certain conditions, the ground state of the system loses this symmetry. Besides the ground state of the system may not be unique but may consist of many degenerate states. One is free to choose any one of the degenerate states as the ground state and express the Lagrangian in terms of new fields which have zero expectation values at the ground state. The new Lagrangian, so obtained, does not possess the same symmetry of the original Lagrangian but physics-wise they are one and the same. This is

[^75]known as the spontaneous symmetry breaking ${ }^{3}$. Let us illustrate it with physical examples choosing (a) real scalar field and (b) complex scalar field.

### 13.3.1 Real scalar field

Consider the Lagrangian of the real scalar field.

$$
\begin{equation*}
\mathscr{L}=T-V=\frac{1}{2}\left(\partial_{\mu} \phi\right)^{2}-\left(\frac{1}{2} \mu^{2} \phi^{2}+\frac{1}{4} \lambda \phi^{4}\right) . \tag{13.25}
\end{equation*}
$$

The Lagrangian is invariant under the symmetry transformation $\phi \rightarrow$ $-\phi$. If both $\lambda$ and $\mu^{2}$ are positive, then the ground state of the system corresponds to $\phi=0$ and the Lagrangian curve will exhibit a reflection symmetry about $\phi=0$, as shown in Fig. 13.1(a). If $\lambda$ is negative, the Lagrangian will go on decreasing with increase of $\phi$ and will not have any minimum value. If $\lambda$ is positive and $\mu^{2}$ is negative, then the Lagrangian will show a minimum at two values of $\phi(\phi= \pm v)$ as shown in Fig. 13.1(b) and they correspond to the ground states of the system.


Figure 13.1: The shape of $V(\phi)$ for a scalar field with $\lambda>0$ and (a) $\mu^{2}>0$ and (b) $\mu^{2}<0$.

$$
\begin{align*}
V & =\frac{1}{2} \mu^{2} \phi^{2}+\frac{1}{4} \lambda \phi^{4} \\
\frac{d V}{d \phi} & =\phi\left(\mu^{2}+\lambda \phi^{2}\right)=0 \tag{13.26}
\end{align*}
$$

[^76]which yields
\[

$$
\begin{equation*}
\phi= \pm \sqrt{-\mu^{2} / \lambda}= \pm v ; \quad \mu^{2}=-\lambda v^{2} . \tag{13.27}
\end{equation*}
$$

\]

In quantum field theory, all perturbation calculations are done with respect to the ground state, $\phi=v$ or $\phi=-v$. Let us choose the ground state to correspond to $\phi=v$. There is no loss of generality since the ground state $\phi=-v$ can always be reached by reflection symmetry.

$$
\begin{equation*}
\phi(\mathbf{x})=v+\eta(\mathbf{x}) \tag{13.28}
\end{equation*}
$$

where $\eta(\mathbf{x})$ is the new field with which all perturbation calculations are done. Let us rewrite the Lagrangian (13.25) in terms of the new field $\eta(\mathbf{x})$.

$$
\begin{align*}
\mathscr{L} & =\frac{1}{2}\left(\partial_{\mu} \eta(\mathbf{x})\right)^{2}-\frac{1}{2} \mu^{2}(v+\eta(\mathbf{x}))^{2}-\frac{1}{4} \lambda(v+\eta(\mathbf{x}))^{4} \\
& =\frac{1}{2}\left(\partial_{\mu} \eta(\mathbf{x})\right)^{2}+\frac{1}{2} \lambda v^{2}(v+\eta(\mathbf{x}))^{2}-\frac{1}{4} \lambda(v+\eta(\mathbf{x}))^{4} \\
& =\frac{1}{2}\left(\partial_{\mu} \eta(\mathbf{x})\right)^{2}+(v+\eta(\mathbf{x}))^{2}\left\{\frac{1}{2} \lambda v^{2}-\frac{1}{4} \lambda(v+\eta(\mathbf{x}))^{2}\right\} . \tag{13.29}
\end{align*}
$$

The quantity within the curly bracket can be simplified to yield

$$
\{\cdots\}=\frac{1}{4} \lambda\left(v^{2}-\eta(\mathbf{x})^{2}-2 v \eta(\mathbf{x})\right) .
$$

Substituting this and rearranging, we get

$$
\begin{equation*}
\mathscr{L}=\frac{1}{2}\left(\partial_{\mu} \eta(\mathbf{x})\right)^{2}-\lambda v^{2} \eta(\mathbf{x})^{2}-\lambda v \eta(\mathbf{x})^{3}-\frac{1}{4} \lambda \eta(\mathbf{x})^{4}+\frac{1}{4} \lambda v^{4} . \tag{13.30}
\end{equation*}
$$

The last term $\frac{1}{4} \lambda v^{4}$ is just a constant. The first term is the kinetic energy term for the $\eta$ field. The second term is the mass term for the scalar field $\eta$, which can be inferred by a reference to Eq. (13.25).

$$
\begin{equation*}
m_{\eta}=\sqrt{2 \lambda v^{2}}=\sqrt{-2 \mu^{2}} . \tag{13.31}
\end{equation*}
$$

The higher order terms in $\eta$ represent the self interaction of the $\eta$ field.
The Lagrangian $\mathscr{L}$ in $\phi$ given by Eq. (13.25) has the reflection symmetry whereas the Lagrangian $\mathscr{L}$ in $\eta$ has no such reflection symmetry. This is known as the spontaneous symmetry breaking. Both the Lagrangians (13.25) and (13.30) are equivalent and they will yield the same physics but the Lagrangian (13.30) which corresponds to the ground state of the system does not possess the symmetry of the Lagrangian (13.25).

### 13.3.2 Complex scalar field

## The Goldstone theorem

Let us now consider the effect of spontaneous symmetry breaking in the case of a complex scalar field $\phi=\left(\phi_{1}+i \phi_{2}\right) / \sqrt{2}$, where $\phi_{1}$ and $\phi_{2}$ are two real fields. The Lagrangian for the complex scalar field can be written as

$$
\begin{equation*}
\mathscr{L}=\left(\partial_{\mu} \phi\right)^{*}\left(\partial_{\mu} \phi\right)-\mu^{2} \phi^{*} \phi-\lambda\left(\phi^{*} \phi\right)^{2} . \tag{13.32}
\end{equation*}
$$

This Lagrangian is invariant under $U(1)$ global gauge transformation.

$$
\phi \rightarrow e^{i \alpha} \phi
$$

As before, we shall consider the interesting case of $\lambda>0$ and $\mu^{2}<0$. The potential

$$
\begin{equation*}
V(\phi)=\mu^{2} \phi^{*} \phi+\lambda\left(\phi^{*} \phi\right)^{2} \tag{13.33}
\end{equation*}
$$

is minimum at the value of $\phi$ for which

$$
\begin{equation*}
\frac{d V}{d \phi}=\mu^{2} \phi^{*}+2 \lambda \phi^{*}\left(\phi^{*} \phi\right)=0 \quad \text { or } \quad|\phi|^{2}=-\frac{\mu^{2}}{2 \lambda}=\frac{v^{2}}{2} \tag{13.34}
\end{equation*}
$$

This corresponds to the vacuum expectation value

$$
\begin{equation*}
\langle\phi\rangle=\frac{v}{\sqrt{2}}=\sqrt{-\frac{\mu^{2}}{2 \lambda}} . \tag{13.35}
\end{equation*}
$$

Only the magnitude is obtained but its phase is arbitrary. There are many ground states located on a circle of radius $v$ with the same energy as shown in Fig. 13.2(b) and the system can be in anyone of the ground states.

We can now expand the complex scalar field $\phi(x)$ in terms of two real fields $\eta(x)$ and $\varepsilon(x)$ which have zero expectation values at the ground state.

$$
\begin{equation*}
\phi(x)=\sqrt{\frac{1}{2}}[v+\eta(x)+i \varepsilon(x)] . \tag{13.36}
\end{equation*}
$$

This yields

$$
\begin{align*}
\phi^{*} \phi & =\frac{1}{2}\left\{(v+\eta)^{2}+\varepsilon^{2}\right\} .  \tag{13.37}\\
V & =\mu^{2}\left(\phi^{*} \phi\right)+\lambda\left(\phi^{*} \phi\right)^{2}, \quad \text { with } \quad \mu^{2}=-v^{2} \lambda \\
& =\frac{\lambda}{4}\left\{(v+\eta)^{2}+\varepsilon^{2}\right\}\left\{-2 v^{2}+(v+\eta)^{2}+\varepsilon^{2}\right\} \\
& =\frac{\lambda}{4}\left\{\left(\eta^{2}+\varepsilon^{2}+2 v \eta\right)^{2}-v^{4}\right\} . \tag{13.38}
\end{align*}
$$



Figure 13.2: The shape of $V(\phi)$ for a complex scalar field with $\lambda>0$ and (a) $\mu^{2}>0$ and (b) $\mu^{2}<0$.

Substituting (13.38) in the Lagrangian density (13.32) and expressing it in terms of the new quantum fields $\eta$ and $\varepsilon$, we get

$$
\begin{align*}
\mathscr{L}= & \frac{1}{2}\left(\partial^{\mu} \eta\right)\left(\partial_{\mu} \eta\right)+\frac{1}{2}\left(\partial^{\mu} \varepsilon\right)\left(\partial_{\mu} \varepsilon\right)-\lambda v^{2} \eta^{2}-\lambda v \eta\left(\eta^{2}+\varepsilon^{2}\right) \\
& -\frac{\lambda}{4}\left(\eta^{2}+\varepsilon^{2}\right)^{2}+\frac{\lambda v^{4}}{4} . \tag{13.39}
\end{align*}
$$

The first two terms are the kinetic energy terms, the third term is the mass term for the $\eta(x)$ field with mass $\sqrt{2 \lambda v^{2}}$, the fourth and fifth terms are the cubic and quartic terms in $\eta$ and $\varepsilon$ and the last term is a constant term which has no physical significance.

We have started with a massive complex scalar field and by spontaneous symmetry breaking, we have obtained a massive scalar field and a massless scalar field with interaction terms. The massive scalar field is known as the gauge boson with mass and the massless scalar field is known as the Goldstone boson. The present study is a simple example of the Goldstone theorem which states that massless scalar bosons appear whenever a continuous symmetry of a physical system is spontaneously broken.

### 13.4 The Higgs mechanism

We have seen earlier that the invariance of the Lagrangian under local phase transformation dictates the interaction dynamics and necessitates the introduction of a gauge field. This gauge field is massless and so
the $U(1)$ gauge transformation describes well the e.m. interaction and one can identify the newly introduced gauge field as the photon. But, in the case of weak interaction, the gauge fields are the intermediate vector Bosons which are massive. In the Lagrangian, the introduction of any term involving the mass of the gauge field will violate the local phase transformation and the theory will become non-renormalizable. So, in order to preserve the local phase transformation symmetry of the original Lagrangian and, at the same time, endow the gauge field with mass, the concept of spontaneous symmetry breaking is used and this method is known as the Higgs mechanism.

The mechanism of mass generation by spontaneous symmetry breaking was proposed in 1964 by three different groups - Englert and Brout ${ }^{4}$, Peter W. Higgs ${ }^{5}$ and Guralink, Hagen and Kibble ${ }^{6}$ - almost simultaneously but working independently. More often, this mechanism is referred to as the Higgs mechanism after one of the discoverers.

### 13.4.1 Interaction of charged scalar field with e.m. field

Let us consider the QED Lagrangian for a charged scalar particle of mass $\mu$, which is invariant under a $\mathrm{U}(1)$ local gauge transformation.

$$
\phi \rightarrow e^{i \alpha(x)} \phi
$$

To obtain such a Lagrangian, the derivative $\partial_{\mu}$ has to be replaced by the covariant derivative $D_{\mu}$

$$
\partial_{\mu} \rightarrow D_{\mu}=\partial_{\mu}-i e A_{\mu},
$$

and the gauge field $A_{\mu}$ transforms as

$$
A_{\mu} \rightarrow A_{\mu}+\frac{1}{e} \partial_{\mu} \alpha
$$

By this procedure, we obtain the gauge-invariant QED Lagrangian for a charged scalar particle.

$$
\begin{equation*}
\mathscr{L}=\left(\partial^{\mu}+i e A^{\mu}\right) \phi^{*}\left(\partial_{\mu}-i e A_{\mu}\right) \phi-V-\frac{1}{4} F_{\mu \nu} F^{\mu \nu}, \tag{13.40}
\end{equation*}
$$

[^77]with
$$
V=\mu^{2} \phi^{*} \phi+\lambda\left(\phi^{*} \phi\right)^{2} .
$$

The Lagrangian (13.40) corresponds to a scalar field $\phi$ of mass $\mu$ and a gauge field $A_{\mu}$ which is massless. By treating $\mu^{2}$ and $\lambda$ as parameters and taking $\mu^{2}<0$ and $\lambda>0$, we will be able to generate a mass for the gauge boson $A^{\mu}$ by the method of spontaneous symmetry breaking. This way of generating a mass for gauge boson is known as the Higgs mechanism.

Minimizing the potential

$$
V=\mu^{2} \phi^{*} \phi+\lambda\left(\phi^{*} \phi\right)^{2}, \quad \frac{d V}{d \phi}=\mu^{2} \phi^{*}+2 \lambda\left(\phi^{*} \phi\right) \phi^{*}=0
$$

we get

$$
|\phi|^{2}=-\frac{\mu^{2}}{2 \lambda}=\frac{v^{2}}{2}, \quad v=\sqrt{-\frac{\mu^{2}}{\lambda}}
$$

Expanding the complex scalar field $\phi$ in terms of the real fields $\eta$ and $\varepsilon$ that correspond to the ground state of the system

$$
\phi(x)=\sqrt{\frac{1}{2}}\{v+\eta(x)+i \varepsilon(x)\},
$$

we can rewrite the Lagrangian $\mathscr{L}$ in terms of the new fields $\eta(x)$ and $\varepsilon(x)$.
The first term in the Lagrangian (13.40) yields

$$
\frac{1}{2}\left(\partial^{\mu} \eta\right)\left(\partial_{\mu} \eta\right)+\frac{1}{2}\left(\partial^{\mu} \varepsilon\right)\left(\partial_{\mu} \varepsilon\right)+e^{2} A^{\mu} A_{\mu} \phi^{*} \phi
$$

Using the expressions (13.37) and (13.38), deduced earlier for $\phi^{*} \phi$ and $V$, we rewrite the Lagrangian (13.40) in terms of the new fields $\eta$ and $\varepsilon$.

$$
\begin{align*}
\mathscr{L}^{\prime}= & \frac{1}{2}\left(\partial^{\mu} \eta\right)\left(\partial_{\mu} \eta\right)+\frac{1}{2}\left(\partial^{\mu} \varepsilon\right)\left(\partial_{\mu} \varepsilon\right)+\frac{1}{2} e^{2} v^{2} A^{\mu} A_{\mu} \\
& -v^{2} \lambda \eta^{2}-e v A_{\mu} \partial^{\mu} \varepsilon-\frac{1}{4} F_{\mu \nu} F^{\mu \nu}+\text { interaction terms. } \tag{13.41}
\end{align*}
$$

The Lagrangian $\mathscr{L}^{\prime}$ exhibits a particle spectrum which consists of a massive gauge vector boson $A_{\mu}$, a massive scalar $\eta$, a massless Goldstone boson $\varepsilon$. From an inspection of the Lagrangian (13.41), we can write down the masses of the particles.

$$
m_{A}=e v, \quad m_{\eta}=\sqrt{2 \lambda v^{2}}, \quad m_{\varepsilon}=0
$$

In this way, we have dynamically generated a mass for the gauge boson $A_{\mu}$ but along with it, we have also a massless Goldstone boson $\varepsilon$. The

Lagrangians $\mathscr{L}$ and $\mathscr{L}^{\prime}$ are equivalent and the number of degrees of freedom cannot change by this transformation. In $\mathscr{L}$, we have a complex scalar field with mass which accounts for two degrees of freedom and a massless gauge boson $A_{\mu}$ with two degrees of freedom corresponding to transverse polarizations. Thus, in total, we have four degrees of freedom in the Lagrangian $\mathscr{L}$. On the other hand, in $\mathscr{L}^{\prime}$, the gauge boson $A_{\mu}$ has acquired a mass and so it can have longitudinal polarization in addition to two transverse polarizations. There are also two scalar particles, one with mass $(\eta)$ and another without mass $(\varepsilon)$ (known as Goldstone boson). Thus, we have, in total, five degrees of freedom in the Lagrangian $\mathscr{L}^{\prime}$. The Goldstone boson, which is spurious, has to be eliminated. The presence of off-diagonal term ev $A_{\mu} \partial^{\mu} \varepsilon$ in $\mathscr{L}^{\prime}$ requires some attention. Is it possible to eliminate the field $\varepsilon$ by choosing some particular gauge transformation? Yes, it is, by choosing a slightly modified gauge transformation.

## The unitary gauge

We get a clue from the expansion of the field $\phi$ in terms of the fields $\eta$ and $\varepsilon$ which define the ground state.

$$
\phi=\sqrt{\frac{1}{2}}(v+\eta+i \varepsilon) .
$$

This is the lowest order in $\varepsilon$ but instead of this, we shall include higher orders in $\varepsilon$ and write

$$
\phi=\sqrt{\frac{1}{2}}(v+\eta) e^{i \varepsilon / v}
$$

This suggests that we should substitute a different set of real fields $h, \theta, A_{\mu}$.

$$
\begin{align*}
\phi \longrightarrow \phi^{\prime} & =\sqrt{\frac{1}{2}}(v+h(x)) e^{i \theta(x) / v}  \tag{13.42}\\
A_{\mu} \longrightarrow A_{\mu}^{\prime} & =A_{\mu}+\frac{1}{e v} \partial_{\mu} \theta \tag{13.43}
\end{align*}
$$

This is a particular choice of gauge, known as the unitary gauge, that makes the Lagrangian independent of the field $\theta$.

$$
\begin{equation*}
\mathscr{L}^{\prime \prime}=\left(\partial^{\mu}+i e A^{\prime \mu}\right) \phi^{\prime *}\left(\partial_{\mu}-i e A_{\mu}^{\prime}\right) \phi^{\prime}-V-\frac{1}{4} F_{\mu \nu} F^{\mu \nu} \tag{13.44}
\end{equation*}
$$

where

$$
\begin{align*}
V & =\mu^{2} \phi^{\prime *} \phi^{\prime}+\lambda\left(\phi^{\prime *} \phi^{\prime}\right)^{2}, \quad \text { with } \quad \mu^{2}=-v^{2} \lambda \\
& =\frac{\lambda}{4}\left(h^{4}+4 v h^{3}+4 v^{2} h^{2}-v^{4}\right) . \tag{13.45}
\end{align*}
$$

The first term in (13.44) yields

$$
\begin{aligned}
\left(\partial^{\mu}\right. & \left.+i e A^{\prime \mu}\right) \phi^{\prime *}\left(\partial_{\mu}-i e A_{\mu}^{\prime}\right) \phi^{\prime} \\
& =\frac{1}{2} \partial^{\mu} h \partial_{\mu} h+\frac{1}{2} e^{2} v^{2} A^{\mu} A_{\mu}+\frac{1}{2} e^{2} h^{2} A^{\mu} A_{\mu}+e^{2} v h A^{\mu} A_{\mu} \cdot(13.46)
\end{aligned}
$$

Substituting (13.45) and (13.46) into the Lagrangian (13.44) and rearranging, we get

$$
\begin{align*}
\mathscr{L}^{\prime \prime}= & \frac{1}{2}\left(\partial_{\mu} h\right)^{2}-\lambda v^{2} h^{2}-\frac{1}{2} e^{2} v^{2} A_{\mu}^{2}-\lambda v h^{3}-\frac{1}{4} \lambda h^{4} \\
& +\frac{1}{2} e^{2} h^{2} A_{\mu}^{2}+v e^{2} h A_{\mu}^{2}+\frac{1}{4} \lambda v^{4}-\frac{1}{4} F^{\mu \nu} F_{\mu \nu} . \tag{13.47}
\end{align*}
$$

We have successfully eliminated the Goldstone boson which is spurious. The Lagrangian (13.47) just describes two interacting massive particles, the massive gauge vector boson $A_{\mu}$ and a massive scalar boson $h$ which is called the Higgs boson.

$$
m_{A}=e v, \quad m_{h}=\sqrt{2 \lambda v^{2}} .
$$

The unwanted massless Goldstone boson has been used up to provide the longitudinal polarization for the massive gauss vector boson $A_{\mu}$. Since the above study offers a method of producing massive gauge vector boson by eliminating the spurious Goldstone boson, this method is also known as the Higgs mechanism.

### 13.4.2 The $\mathrm{SU}(2)$ gauge symmetry

Hitherto, we have considered the $\mathrm{U}(1)$ gauge symmetry of the Lagrangian and its spontaneous symmetry breaking. The $\mathrm{U}(1)$ gauge group is abelian and so it is much simpler to deal with. Here, we shall consider the $\mathrm{SU}(2)$ gauge symmetry which is a bit more complicated than the $U(1)$ gauge group because the $\mathrm{SU}(2)$ gauge group is non-abelian. Non-abelian gauge theories are often referred to as Yang-Mills theories, named after the persons who have first discovered their structure.

The generators of the $\mathrm{SU}(2)$ group are the Pauli spin matrices $\tau_{1}, \tau_{2}, \tau_{3}$ which do not commute.

$$
\begin{gather*}
\tau_{1}=\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right], \quad \tau_{2}=\left[\begin{array}{rr}
0 & -i \\
i & 0
\end{array}\right] \quad \tau_{3}=\left[\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right] .  \tag{13.48}\\
{\left[\tau_{i}, \tau_{j}\right]_{-}=2 i \epsilon_{i j k} \tau_{k} .} \tag{13.49}
\end{gather*}
$$

In Eq. (13.49), $\epsilon_{i j k}$ denotes an antisymmetric tensor. $\epsilon_{i j k}=1$, if $i, j, k$ are in cyclic order; otherwise $\epsilon_{i j k}=-1$.

Let us consider a Lagrangian

$$
\begin{equation*}
\mathscr{L}=\left(\partial_{\mu} \phi\right)^{\dagger}\left(\partial^{\mu} \phi\right)-\mu^{2} \phi^{\dagger} \phi-\lambda\left(\phi^{\dagger} \phi\right)^{2}, \tag{13.50}
\end{equation*}
$$

where $\phi$ is an $\mathrm{SU}(2)$ doublet of complex scalar fields.

$$
\phi=\left[\begin{array}{c}
\phi_{\alpha}  \tag{13.51}\\
\phi_{\beta}
\end{array}\right]=\sqrt{\frac{1}{2}}\left[\begin{array}{c}
\phi_{1}+i \phi_{2} \\
\phi_{3}+i \phi_{4}
\end{array}\right] .
$$

It can be easily checked that the Lagrangian $\mathscr{L}$ is invariant under global $\mathrm{SU}(2)$ phase transformation.

$$
\begin{equation*}
\phi \rightarrow \phi^{\prime}=e^{i \alpha_{a} \tau_{a} / 2} \phi . \tag{13.52}
\end{equation*}
$$

To make it invariant under local gauge transformation with $\alpha(\mathbf{x})$, we need to replace the derivative $\partial_{\mu}$ with covariant derivative

$$
\begin{equation*}
D_{\mu}=\partial_{\mu}+\frac{1}{2} i g \tau_{a} W_{\mu}^{a} \tag{13.53}
\end{equation*}
$$

invoking three gauge vector fields $W_{\mu}^{a}$ with $a=1,2,3$. The strength of the $\mathrm{SU}(2)$ coupling to the Gauge fields is denoted by $g$. Under an infinitesimal gauge transformation

$$
\begin{equation*}
\phi(\mathbf{x}) \rightarrow \phi^{\prime}(\mathbf{x})=(1+i \boldsymbol{\alpha}(\mathbf{x}) \cdot \boldsymbol{\tau} / 2) \phi(\mathbf{x}) \tag{13.54}
\end{equation*}
$$

the three gauge fields transform as (vide Solved Problems)

$$
\begin{equation*}
\boldsymbol{W}_{\mu} \rightarrow \boldsymbol{W}_{\mu}-\frac{1}{g} \partial_{\mu} \boldsymbol{\alpha}-\boldsymbol{\alpha} \times \boldsymbol{W}_{\mu} . \tag{13.55}
\end{equation*}
$$

The extra term $\boldsymbol{\alpha} \times \boldsymbol{W}_{\mu}$ in Eq. (13.55) occurs because $\boldsymbol{W}_{\mu}$ is a $\mathrm{SU}(2)$ vector. With these inputs, we can write the gauge invariant Lagrangian $\mathscr{L}$ as

$$
\begin{align*}
\mathscr{L}= & \left(\partial_{\mu} \phi+\frac{1}{2} i g \boldsymbol{\tau} \cdot \boldsymbol{W}_{\mu} \phi\right)^{\dagger}\left(\partial^{\mu} \phi+\frac{1}{2} i g \boldsymbol{\tau} \cdot \boldsymbol{W}^{\mu} \phi\right) \\
& -V(\phi)-\frac{1}{4} \boldsymbol{W}_{\mu \nu} \cdot \boldsymbol{W}^{\mu \nu} \tag{13.56}
\end{align*}
$$

with

$$
\begin{align*}
V(\phi) & =\mu^{2} \phi^{\dagger} \phi+\lambda\left(\phi^{\dagger} \phi\right)^{2}  \tag{13.57}\\
\boldsymbol{W}_{\mu \nu} & =\partial_{\mu} \boldsymbol{W}_{\nu}-\partial_{\nu} \boldsymbol{W}_{\mu}-g \boldsymbol{W}_{\mu} \times \boldsymbol{W}_{\nu} \tag{13.58}
\end{align*}
$$

The extra terms in Eqs. (13.55) and (13.58) arise due to the non-abelian character of the $\mathrm{SU}(2)$ group since $\tau_{a}, a=1,2,3$ do not commute with each other. If $\mu^{2}>0$ and $\lambda>0$, then the Lagrangian (13.56) describes a system of four scalar particles, each of mass $\mu$, interacting with three massless vector bosons.

What we are interested in, is the special case of $\mu^{2}<0$ and $\lambda>0$, when a spontaneous symmetry breaking occurs. Out of the four scalar particles, three of them are absorbed by the three gauge vector bosons to gain mass and the remaining one becomes the Higgs boson. This method of adding mass to the Gauge bosons is known as the Higgs mechanism.

The potential $V(\phi)$ has its minimum when

$$
\begin{equation*}
\phi^{\dagger} \phi=\frac{1}{2}\left(\phi_{1}^{2}+\phi_{2}^{2}+\phi_{3}^{2}+\phi_{4}^{2}\right)=-\frac{\mu^{2}}{2 \lambda} . \tag{13.59}
\end{equation*}
$$

There are a large number of points at which $V(\phi)$ is minimum and one is free to choose anyone of them. Let us choose the point

$$
\begin{equation*}
\phi_{1}=\phi_{2}=\phi_{4}=0, \quad \phi_{3}^{2}=-\frac{\mu^{2}}{\lambda}=v^{2} \tag{13.60}
\end{equation*}
$$

as the minimum and expand $\phi(\mathbf{x})$ about this point. This is what is known as the spontaneous symmetry breaking of the $\mathrm{SU}(2)$ symmetry. The symmetry which is explicit in (13.59) is broken in (13.60). Substituting (13.60) in Eq. (13.51), we get the vacuum state $\phi_{0}$.

$$
\phi_{0}=\sqrt{\frac{1}{2}}\left[\begin{array}{l}
0 \\
v
\end{array}\right]
$$

One can expand $\phi(\mathbf{x})$ about this vacuum state, postulating a Higgs field $h(\mathbf{x})$.

$$
\phi(\mathbf{x})=\sqrt{\frac{1}{2}}\left[\begin{array}{c}
0  \tag{13.61}\\
v+h(\mathbf{x})
\end{array}\right] .
$$

By substituting (13.61) for $\phi(\mathbf{x})$ in the Lagrangian (13.56), we find that the Lagrangian now depends on only one field, called the Higgs field. It is really surprising. We initially started with four scalar fields $\phi_{1}, \phi_{2}, \phi_{3}, \phi_{4}$ but ended up with only one field $h$. Let us consider the small $\mathrm{SU}(2)$ fluctuations about the vacuum in terms of the real fields $\theta_{1}, \theta_{2}, \theta_{3}$ and $h$.

$$
\phi(\mathbf{x})=e^{i \boldsymbol{\tau} \cdot \boldsymbol{\theta}(\mathbf{x}) / v}\left[\begin{array}{c}
0  \tag{13.62}\\
\frac{v+h(\mathbf{X})}{\sqrt{2}}
\end{array}\right] .
$$

For small perturbations,

$$
\begin{align*}
e^{i \boldsymbol{\tau} \cdot \boldsymbol{\theta}(\mathbf{x}) / v} & =1+\frac{i}{v}\left(\tau_{1} \theta_{1}+\tau_{2} \theta_{2}+\tau_{3} \theta_{3}\right) \\
& =\left[\begin{array}{cc}
1+i \theta_{3} / v & \left(\theta_{2}+i \theta_{1}\right) / v \\
-\left(\theta_{2}-i \theta_{1}\right) / v & 1-i \theta_{3} / v
\end{array}\right] \tag{13.63}
\end{align*}
$$

Substituting (13.63) in (13.62), we get

$$
\phi(\mathbf{x}) \approx \sqrt{\frac{1}{2}}\left[\begin{array}{c}
\theta_{2}+i \theta_{1}  \tag{13.64}\\
v+h-i \theta_{3}
\end{array}\right]
$$

where we have retained only the first order terms, neglecting the secondorder terms. Thus, we find that $\phi(\mathbf{x})$ indeed consists of four independent fields and parametrize the fluctuations from the vacuum. Since the Lagrangian is locally $\mathrm{SU}(2)$ invariant, we can gauge away the three massless Goldstone bosons $\theta_{1}(\mathbf{x}), \theta_{2}(\mathbf{x}), \theta_{3}(\mathbf{x})$, giving mass to the three gauge bosons $\boldsymbol{W}$.

The last term in the Lagrangian (13.56) represents the kinetic energy of the gauge bosons. Let us substitute $\phi(\mathbf{x})$ given by Eq. (13.61) in the remaining terms of the Lagrangian (13.56) and consider the relevant non-vanishing parts.

$$
\begin{align*}
\left(\partial_{\mu} \phi\right)^{\dagger}\left(\partial^{\mu} \phi\right) & =\frac{1}{2}\left(\partial_{\mu} h\right)^{\dagger}\left(\partial^{\mu} h\right) .  \tag{13.65}\\
\left|\frac{1}{2} i g \boldsymbol{\tau} \cdot \boldsymbol{W}_{\mu} \phi\right|^{2} & =\frac{g^{2}}{8}\left|\left[\begin{array}{cc}
W_{\mu}^{3} & W_{\mu}^{1}-i W_{\mu}^{2} \\
W_{\mu}^{1}+i W_{\mu}^{2} & W_{\mu}^{3}
\end{array}\right]\left[\begin{array}{c}
0 \\
v+h
\end{array}\right]\right|^{2} \\
& =\frac{g^{2}}{8}\left[\left|W_{\mu}^{1}\right|^{2}+\left|W_{\mu}^{2}\right|^{2}+\left|W_{\mu}^{3}\right|^{2}\right](v+h)^{2} .  \tag{13.66}\\
\phi^{\dagger} \phi & =\frac{1}{2}(v+h)^{2} .  \tag{13.67}\\
V(\phi) & =\phi^{\dagger} \phi\left\{\mu^{2}+\lambda\left(\phi^{\dagger} \phi\right)\right\} \\
& =\frac{1}{2}(v+h)^{2}\left\{-v^{2} \lambda+\frac{\lambda}{2}(v+h)^{2}\right\} \\
& =\frac{\lambda}{4}\left\{\left(h^{2}+2 v h\right)^{2}-v^{4}\right\} . \tag{13.68}
\end{align*}
$$

Equation (13.65) represents the kinetic energy of the Higgs boson. In Eq. (13.66), we have used the short-hand notation $\left|\left.\right|^{2}\right.$ to denote ( ) ${ }^{\dagger}()$ and it yields the mass of the gauge bosons $W_{\mu}^{a}$ besides their interaction with the Higgs field. From the term

$$
\left(g^{2} v^{2} / 8\right)\left[\left|W_{\mu}^{1}\right|^{2}+\left|W_{\mu}^{2}\right|^{2}+\left|W_{\mu}^{3}\right|^{2}\right]
$$

one can infer that all the three gauge bosons have identical masses $M=$ $\frac{1}{2} g v$ since the typical mass term is of the form $\frac{1}{2} M^{2} W_{\mu}^{2}$. From Eq. (13.68), we also infer that the mass of the Higgs boson is $m=\sqrt{2 \lambda v^{2}}$ from the typical term $\frac{1}{2} m^{2} h^{2}=\lambda v^{2} h^{2}$.

$$
M=\frac{1}{2} g v, \quad m=\sqrt{2 \lambda v^{2}} .
$$

Thus, we find that the Lagrangian (13.56) describes three gauge bosons, each with mass $\frac{1}{2} g v$ and a massive Higgs boson with mass $\sqrt{2 \lambda v^{2}}$. By using the unitary gauge, the three massless Goldstone bosons that were created by spontaneous symmetry breaking were swallowed by the three gauge bosons in order to gain mass. The scalar degrees of freedom have been used to provide the longitudinal polarization of the massive gauge bosons. This is an illustrative example of the Higgs mechanism.

### 13.5 Electroweak Interactions

Having learnt how a Lagrangian, which is invariant under global phase transformation, can be made invariant under local phase transformation by introducing vector gauge bosons and having studied the concepts of spontaneous symmetry breaking and the Higgs mechanism, we are now in a position to discuss the gauge theory of electroweak interaction.

Since the neutrino is left-handed, it is coupled with the left-handed electron to form a weak isospin doublet. The right-handed electron is considered as an isospin singlet. In this scheme, the Dirac Lagrangian of the electron-neutrino pair is invariant under global $\mathrm{SU}(2) \times \mathrm{U}(1)$ gauge transformation, only if the mass term of the electron is excluded from the Lagrangian. In order to make the Lagrangian invariant under the local $\mathrm{SU}(2) \times \mathrm{U}(1)$ gauge transformation, we need to introduce four vector gauge bosons $\boldsymbol{W}_{1}, \boldsymbol{W}_{2}, \boldsymbol{W}_{3}$ and $\boldsymbol{B}$ and by introducing scalar fields (Higgs fields) and invoking the Higgs mechanism, these vector gauge bosons are made to acquire masses. A suitable linear combinations of these gauge bosons yield the three massive vector bosons $\boldsymbol{W}^{+}, \boldsymbol{W}^{-}$and $\boldsymbol{Z}^{0}$, associated with weak interaction and the massless photon $\boldsymbol{A}$, associated with electromagnetic interaction. The same Higgs doublet of scalar fields is used to generate lepton masses by spontaneous symmetry breaking but a slightly modified Higgs doublet of scalar fields is used to generate quark masses. This is the gauge theory of electroweak interaction that has been developed by Abdus Salam and Weinberg.

### 13.5.1 Weak isospin and weak hypercharge

Just as the strongly interacting particles (hadrons) are grouped together in an elegant way in the Gellmann-Nishijima scheme by defining isospin and hypercharge quantum numbers, the weakly interacting particles are also grouped together in an analogous way by defining weak isospin $T$ and weak hypercharge $Y$ quantum numbers as shown in Table 13.1 such that the charge of the particle is given by

$$
\begin{equation*}
Q=T_{3}+\frac{Y}{2} \tag{13.69}
\end{equation*}
$$

There are three generations of leptons and quarks.

$$
\binom{\nu_{e}}{e^{-}},\binom{u}{d} ; \quad\binom{\nu_{\mu}}{\mu^{-}}, \quad\binom{c}{s} ; \quad\binom{\nu_{\tau}}{\tau^{-}}, \quad\binom{t}{b}
$$

In Table 13.1, the weak isospin $T$ and weak hypercharge $Y$ quantum numbers are given only for the firat generation of leptons and quarks. Similar quantum numbers are given for the second and third generations of leptons and quarks.

Table 13.1: Weak isospin and weak hypercharge quantum numbers of the first generation of leptons and quarks

| Lepton | $T$ | $T_{3}$ | $Q$ | $Y$ | Quark | $T$ | $T_{3}$ | $Q$ | $Y$ |
| :---: | :---: | ---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\nu_{e}$ | $\frac{1}{2}$ | $\frac{1}{2}$ | 0 | -1 | $u_{L}$ | $\frac{1}{2}$ | $\frac{1}{2}$ | $\frac{2}{3}$ | $\frac{1}{3}$ |
| $e_{L}^{-}$ | $\frac{1}{2}$ | $-\frac{1}{2}$ | -1 | -1 | $d_{L}$ | $\frac{1}{2}$ | $-\frac{1}{2}$ | $-\frac{1}{3}$ | $\frac{1}{3}$ |
|  |  |  |  |  | $u_{R}$ | 0 | 0 | $\frac{2}{3}$ | $\frac{4}{3}$ |
| $e_{R}^{-}$ | 0 | 0 | -1 | -2 | $d_{R}$ | 0 | 0 | $-\frac{1}{3}$ | $-\frac{2}{3}$ |

It is an experimental observation that the neutrino is left-handed and this supports the attribute of zero mass to the neutrino ${ }^{7}$. Since the electron has a finite mass, it can have both right-handed and left-handed components. In other words, electron can exist in both chiral states or

[^78]with positive and negative helicities. As can be observed from the Table 13.1, the electron neutrino and the left chiral electron are paired together to form a doublet with weak iso-spin $T=\frac{1}{2}$ and weak hypercharge $Y=-1$ but different iso-spin projections $T_{3}=+\frac{1}{2}$ and $T_{3}=-\frac{1}{2}$. The right chiral electron is a singlet with weak iso-spin $T=0$ and weak hypercharge $Y=-2$.
\[

$$
\begin{array}{ll}
\psi_{L}=\left[\begin{array}{c}
\nu_{e} \\
e_{L}^{-}
\end{array}\right] & \text {with } T=\frac{1}{2}, Y=-1 \\
\psi_{R}=e_{R}^{-} & \text {with } T=0, Y=-2 . \tag{13.71}
\end{array}
$$
\]

They are subjected to gauge transformation $S U(2) \times U(1)$ as shown below:

$$
\begin{align*}
\psi_{L} & \longrightarrow \psi_{L}^{\prime}=e^{i \boldsymbol{\alpha} \cdot \boldsymbol{T}+i \beta(Y / 2)} \psi_{L} ;  \tag{13.72}\\
\psi_{R} & \longrightarrow \psi_{R}^{\prime}=e^{i \beta(Y / 2)} \psi_{R} . \tag{13.73}
\end{align*}
$$

For global phase invariance, the quantities $\alpha$ and $\beta$ will be independent of space time coordinates and for local phase invariance, $\boldsymbol{\alpha}(\mathbf{x})$ and $\beta(\mathbf{x})$ will be dependent on the space-time coordinates $\mathbf{x}$.

It can be easily checked that the mass term for the electron is not gauge invariant since the left-handed electron $e_{L}$ is a member of weak isospin doublet and the right-handed electron $e_{R}$ is a weak isospin singlet. (For convenience of notation, the fields are sometimes represented by the particle symbol itself.)

$$
\begin{align*}
m \bar{e} e & =m\left(\bar{e}_{L}+\bar{e}_{R}\right)\left(e_{L}+e_{R}\right) \\
& =m\left(\bar{e}_{L} e_{L}+\bar{e}_{L} e_{R}+\bar{e}_{R} e_{L}+\bar{e}_{R} e_{R}\right) . \tag{13.74}
\end{align*}
$$

The cross terms are not manifestly gauge invariant and so the mass term for the fermions cannot be included in the Lagrangian, if it is to be gauge invariant.

### 13.5.2 Lagrangian for the fermions

In the doublet structure proposed for the leptons, the mass term is not gauge invariant. So, omitting the mass term, the Dirac Lagrangian for the electron-neutrino pair can be written as

$$
\mathscr{L}=\bar{\psi}_{L} \gamma^{\mu}\left(i \partial_{\mu}\right) \psi_{L}+\bar{e}_{R} \gamma^{\mu}\left(i \partial_{\mu}\right) e_{R},
$$

which is invariant under global gauge transformation. In order to make it invariant under local gauge transformation, we need to introduce four vector gauge bosons $W_{1}, W_{2}, W_{3}$ and $B$ in order to make the derivative $\partial_{\mu}$ to correspond to a covariant derivative $D_{\mu}$.

$$
\begin{aligned}
\left(i \partial_{\mu}\right)_{L} \rightarrow\left(i D_{\mu}\right)_{L} & =\left(i \partial_{\mu}\right)_{L}-g \boldsymbol{T} \cdot \boldsymbol{W}-g^{\prime} \frac{1}{2} Y B_{\mu}, \text { with } Y
\end{aligned}=-1 ;(13.75), \text { with } Y=-2 .(13.76)
$$

Thus, we obtain the $S U(2) \times U(1)$ local gauge invariant Lagrangian of the electron-neutrino pair, by substituting $T=\frac{1}{2} \tau$ and the respective $Y$ values.

$$
\begin{align*}
\mathscr{L}_{1}= & \bar{\psi}_{L} \gamma^{\mu}\left(i \partial_{\mu}-\frac{1}{2} g \boldsymbol{\tau} \cdot \boldsymbol{W}+\frac{1}{2} g^{\prime} B_{\mu}\right) \psi_{L}+\bar{e}_{R} \gamma^{\mu}\left(i \partial_{\mu}+g^{\prime} B_{\mu}\right) e_{R} \\
& -\frac{1}{4} \boldsymbol{W}_{\mu \nu} \cdot \boldsymbol{W}^{\mu \nu}-\frac{1}{4} B_{\mu \nu} B^{\mu \nu} . \tag{13.77}
\end{align*}
$$

The last two terms in Eq. (13.77) represent the kinetic energies and selfcoupling of the $\boldsymbol{W}_{\mu}$ fields and the kinetic energy of the $B_{\mu}$ field.

$$
\begin{align*}
\boldsymbol{W}_{\mu \nu} & =\partial_{\mu} \boldsymbol{W}_{\nu}-\partial_{\nu} \boldsymbol{W}_{\mu}-g \boldsymbol{W}_{\mu} \times \boldsymbol{W}_{\nu}  \tag{13.78}\\
B_{\mu \nu} & =\partial_{\mu} B_{\nu}-\partial_{\nu} B_{\mu} . \tag{13.79}
\end{align*}
$$

Note that the Lagrangian $\mathscr{L}_{1}$ corresponds to massless fermions and massless gauge bosons. The mass terms such as

$$
m \bar{e} e, \quad \frac{1}{2} M_{W}^{2} \boldsymbol{W}^{\mu} \boldsymbol{W}_{\mu}, \quad \frac{1}{2} M_{B}^{2} B^{\mu} B_{\mu}
$$

are not gauge invariant and inclusion of these in the Lagrangian $\mathscr{L}_{1}$ will spoil the guage invariance of the Lagrangian. So, to generate the masses in a gauge invariant way, we invoke the Higgs mechanism by spontaneously breaking the gauge symmetry. This has the great virtue of preserving the gauge invariance of the original Lagrangian, which is considered so essential for retaining the theory as fully renormalizable.

### 13.5.3 The Higgs field

To activate the Higgs mechanism for generating masses for the gauge bosons, we have to invoke the scalar fields $\phi_{i}$ (the Higgs fields) and consider their $S U(2) \times U(1)$ gauge-invariant Lagrangian $\mathscr{L}_{2}$.

$$
\begin{equation*}
\mathscr{L}_{2}=\left|\left(i \partial_{\mu}-g \boldsymbol{T} \cdot \boldsymbol{W}_{\mu}-g^{\prime} \frac{Y}{2} B_{\mu}\right) \phi\right|^{2}-V(\phi), \tag{13.80}
\end{equation*}
$$

with

$$
V(\phi)=\mu^{2} \phi^{\dagger} \phi+\lambda\left(\phi^{\dagger} \phi\right)^{2}
$$

In Eq. (13.80), $g$ and $g^{\prime}$ denote the weak coupling constants, coupling the gauge bosons with the Higgs fields and | $\left.\right|^{2}=()^{\dagger}()$.

It is known that the electro-weak interaction is generated by three massive gauge bosons $W^{+}, W^{-}$and $Z^{0}$ besides the massless photon. Let us bear this in mind, while choosing the scalar fields.

In order to maintain the gauge invariance of $\mathscr{L}_{2}$, one should choose the scalar fields $\phi_{i}$ in such a way that they form a $S U(2) \times U(1)$ multiplet (an isospin doublet with weak hypercharge $Y=1$ )

$$
\phi=\left[\begin{array}{l}
\phi_{\alpha}  \tag{13.81}\\
\phi_{\beta}
\end{array}\right]
$$

with

$$
\begin{equation*}
\phi_{\alpha}=\frac{1}{\sqrt{2}}\left(\phi_{1}+i \phi_{2}\right) ; \quad \phi_{\beta}=\frac{1}{\sqrt{2}}\left(\phi_{3}+i \phi_{4}\right) . \tag{13.82}
\end{equation*}
$$

This is the choice made originally by Weinberg in 1967 and it is popularly known as the Weinberg-Salam model.

Let us minimize the Higgs potential $V(\phi)$ by choosing $\mu^{2}<0$ and $\lambda>0$ and obtain the condition

$$
\begin{equation*}
\phi^{\dagger} \phi=\frac{1}{2}\left(\phi_{1}^{2}+\phi_{2}^{2}+\phi_{3}^{2}+\phi_{4}^{2}\right)=-\frac{\mu^{2}}{2 \lambda} . \tag{13.83}
\end{equation*}
$$

There are many sets of values of $\phi_{i}$ at which $V(\phi)$ is minimum and one is free to choose anyone of them. Let us choose the configuration

$$
\begin{equation*}
\phi_{1}=\phi_{2}=\phi_{4}=0, \quad \phi_{3}^{2}=-\frac{\mu^{2}}{\lambda}=v^{2} \tag{13.84}
\end{equation*}
$$

for the minimum of $V(\phi)$ and identify it as the ground state with the vacuum expectation value $\phi_{0}$.

$$
\phi_{0}=\frac{1}{\sqrt{2}}\left[\begin{array}{l}
0  \tag{13.85}\\
v
\end{array}\right]
$$

The choice of this vacuum $\phi_{0}$ spontaneously breaks the $S U(2) \times U(1)_{Y}$ symmetry of the Lagrangian $\mathscr{L}_{2}$ and thereby the gauge bosons acquire masses. But there is a subgroup of gauge transformations whose symmetry is not broken by the vacuum $\phi_{0}$ and so the gauge boson associated
with this subgroup will remain massless. The choice of the vacuum $\phi_{0}$ with $T=\frac{1}{2}, T_{3}=-\frac{1}{2}$ and $Y=1$ preserves the electromagnetic $U(1)_{\mathrm{em}}$ symmetry with the generator

$$
\begin{equation*}
Q=T_{3}+\frac{Y}{2}, \tag{13.86}
\end{equation*}
$$

such that

$$
\begin{equation*}
Q \phi_{0}=0 . \tag{13.87}
\end{equation*}
$$

Consequently,

$$
\begin{equation*}
\phi_{0} \rightarrow \phi_{0}^{\prime}=e^{i \alpha(\mathbf{x}) Q} \phi_{0}=\phi_{0} . \tag{13.88}
\end{equation*}
$$

The gauge invariance of the vacuum under $U(1)_{\mathrm{em}}$ transformation assures that the photon which is the gauge boson of electromagnetic interaction is massless.

The four gauge bosons $\boldsymbol{W}_{1}, \boldsymbol{W}_{2}, \boldsymbol{W}_{3}$ and $\boldsymbol{B}$ acquire masses due to the spontaneous symmetry breaking of the $S U(2) \times U(1)_{Y}$ symmetry. The two vector bosons $\boldsymbol{W}_{1}$ and $\boldsymbol{W}_{2}$ combine to yield the two charged intermediate vector bosons $\boldsymbol{W}^{+}$and $\boldsymbol{W}^{-}$.

$$
\begin{equation*}
\boldsymbol{W}^{+}=\frac{1}{\sqrt{2}}\left(\boldsymbol{W}_{1}+i \boldsymbol{W}_{2}\right) ; \quad \boldsymbol{W}^{-}=\frac{1}{\sqrt{2}}\left(\boldsymbol{W}_{1}-i \boldsymbol{W}_{2}\right) . \tag{13.89}
\end{equation*}
$$

On the other hand, the two neutral gauge bosons $\boldsymbol{W}_{3}$ and $\boldsymbol{B}$ combine to yield the neutral intermediate vector boson $\boldsymbol{Z}^{0}$ and the photon $\boldsymbol{A}$.

$$
\begin{align*}
A_{\mu} & =\cos \theta_{W} B_{\mu}+\sin \theta_{W} W_{\mu}^{3}  \tag{13.90}\\
Z_{\mu}^{0} & =-\sin \theta_{W} B_{\mu}+\cos \theta_{W} W_{\mu}^{3} \tag{13.91}
\end{align*}
$$

where $\theta_{W}$ is known as the Weinberg angle or the weak mixing angle. In other words, the eigenstates $A_{\mu}$ and $Z_{\mu}^{0}$ are obtained from the eigenstates $B_{\mu}$ and $W_{\mu}^{3}$ by a unitary transformation which corresponds to a rotation in the Hilbert space through an angle $\theta_{W}$.

Let us now consider how the gauge bosons acquire masses by spontaneous symmetry breaking by the vacuum state $\phi_{0}$. For this consider the
relevant terms in the Lagrangian $\mathscr{L}_{2}$ (with $Y=1$ ).

$$
\begin{align*}
& \left|\left(-g \frac{1}{2} \boldsymbol{\tau} \cdot \boldsymbol{W}_{\mu}-g^{\prime} \frac{Y}{2} B_{\mu}\right) \phi_{0}\right|^{2} \\
& \quad=\frac{1}{8}\left|\left[\begin{array}{cc}
g W_{\mu}^{3}+g^{\prime} B_{\mu} & g\left(W_{\mu}^{1}-i W_{\mu}^{2}\right) \\
g\left(W_{\mu}^{1}+i W_{\mu}^{2}\right) & -g W_{\mu}^{3}+g^{\prime} B_{\mu}
\end{array}\right]\left[\begin{array}{c}
0 \\
v
\end{array}\right]\right|^{2} \\
& \quad=\frac{1}{8} v^{2}\left\{g^{2}\left|W_{\mu}^{1}-i W_{\mu}^{2}\right|^{2}+\left|-g W_{\mu}^{3}+g^{\prime} B_{\mu}\right|^{2}\right\} \\
& \quad=\frac{1}{8} v^{2}\left\{g^{2}\left(2 W_{\mu}^{+} W^{-\mu}\right)+\left(g^{2} W_{3} \mu W^{3 \mu}-2 g g^{\prime} W_{\mu}^{3} B^{\mu}+g^{\prime 2} B_{\mu} B^{\mu}\right)\right\} \tag{13.92}
\end{align*}
$$

where

$$
\boldsymbol{W}^{ \pm}=\frac{1}{\sqrt{2}}\left(\boldsymbol{W}^{1} \pm i \boldsymbol{W}^{2}\right) .
$$

The first term represents the mass term for the charged bosons, $M_{W}^{2} W^{+} W^{-}$ with

$$
\begin{equation*}
M_{W}=\frac{1}{2} v g . \tag{13.93}
\end{equation*}
$$

The other term is off-diagonal in the $W_{\mu}^{(3)}$ and $B_{\mu}$ basis and can be written as

$$
\begin{align*}
& \frac{1}{8} v^{2}\left\{g^{2} W_{3} \mu W^{3 \mu}-2 g g^{\prime} W_{\mu}^{3} B^{\mu}+g^{\prime 2} B_{\mu} B^{\mu}\right\} \\
& \quad=\frac{1}{8} v^{2}\left\{\left|g W_{\mu}^{3}-g^{\prime} B_{\mu}\right|^{2}+0\left|g^{\prime} W_{\mu}^{3}+g B_{\mu}\right|^{2}\right\} . \tag{13.94}
\end{align*}
$$

Eq. (13.94) suggests that instead of two neutral fields $W_{\mu}^{(3)}$ and $B_{\mu}$, we can take a linear combination of them to correspond to two orthogonal fields, one to represent the neutral intermediate vector boson $Z_{\mu}^{0}$ and the other to represent the photon $A_{\mu}$, with masses $M_{z}$ and $M_{A}$ and identify Eq. (13.94) with

$$
\frac{1}{2} M_{Z}^{2}\left(Z_{\mu}^{0}\right)^{2}+\frac{1}{2} M_{A}^{2} A_{\mu}^{2}
$$

where the factor $\frac{1}{2}$ is the appropriate factor for the mass terms of neutral vector bosons $Z_{\mu}^{0}$ and $A_{\mu}$.

The normalized fields $Z_{\mu}^{0}$ anf $A_{\mu}$ are

$$
\begin{align*}
Z_{\mu}^{0} & =\frac{g W_{\mu}^{(3)}-g^{\prime} B_{\mu}}{\sqrt{g^{2}+g^{\prime 2}}} \quad \text { with } M_{Z}=\frac{1}{2} v \sqrt{g^{2}+g^{\prime 2}}  \tag{13.95}\\
A_{\mu} & =\frac{g^{\prime} W_{\mu}^{(3)}+g B_{\mu}}{\sqrt{g^{2}+g^{\prime 2}}} \quad \text { with } M_{A}=0 \tag{13.96}
\end{align*}
$$

It is conventional to relate the electromagnetic coupling constant $e$ to the weak coupling constants $g$ and $g^{\prime}$ by defining a mixing angle $\theta_{W}$ known as the Weinberg angle.

$$
\begin{equation*}
e=g \sin \theta_{W}=g^{\prime} \cos \theta_{W} \tag{13.97}
\end{equation*}
$$

In terms of the mixing angle $\theta_{W}$, Eqs. (13.95) and (13.96) can be written as

$$
\left[\begin{array}{l}
Z_{\mu}^{0}  \tag{13.98}\\
A_{\mu}
\end{array}\right]=\left[\begin{array}{rr}
\cos \theta_{W} & -\sin \theta_{W} \\
\sin \theta_{W} & \cos \theta_{W}
\end{array}\right]\left[\begin{array}{l}
W_{\mu}^{(3)} \\
B_{\mu}
\end{array}\right]
$$

From Eq. (13.97), it follows that

$$
\begin{equation*}
\frac{g^{\prime}}{g}=\tan \theta_{W} . \tag{13.99}
\end{equation*}
$$

From Eqs. (13.93) and (13.95), we obtain the ratio of the masses $M_{W} / M_{Z}$.

$$
\begin{equation*}
\frac{M_{W}}{M_{Z}}=\frac{1}{\sqrt{1+\left(g^{\prime 2} / g^{2}\right)}}=\frac{1}{\sqrt{1+\tan ^{2} \theta_{W}}}=\cos \theta_{W} \tag{13.100}
\end{equation*}
$$

In summary, it can be said that the photon mass $M_{A}=0$ is an input but the ratio of the masses $M_{W} / M_{Z}$ is a definite prediction of the WeinbergSalam model.

### 13.5.4 Lepton masses

It may be recalled that in the original Lagrangian $\mathscr{L}_{1}$, we have not included the mass term $m \bar{\psi} \psi$ for the lepton since it is not invariant under $S U(2) \times U(1)$ gauge transformation. But it is found that the same Higgs doublet which is used to generate masses for the intermediate vector bosons can also be used to generate the masses for the leptons by spontaneous symmetry breaking. For this purpose, let us include the following $S U(2) \times U(1)$ guage invariant Lagrangian $\mathscr{L}_{3}$.

$$
\mathscr{L}_{3}=-G_{e}\left\{\left[\bar{\nu}_{e}, \bar{e}\right]_{L}\left[\begin{array}{c}
\phi_{\alpha}  \tag{13.101}\\
\phi_{\beta}
\end{array}\right] e_{R}+\bar{e}_{R}\left[\bar{\phi}_{\alpha}, \bar{\phi}_{\beta}\right]\left[\begin{array}{c}
\nu_{e} \\
e
\end{array}\right]_{L}\right\} .
$$

The Higgs doublet has exactly the quantum numbers required to make the terms $\bar{e}_{L} e_{R}$ and $\bar{e}_{R} e_{L}$, gauge invariant. By spontaneous symmetry breaking, we obtain

$$
\phi_{0}=\frac{1}{\sqrt{2}}\left[\begin{array}{l}
0  \tag{13.102}\\
v
\end{array}\right] ; \quad \phi(\mathbf{x})=\frac{1}{\sqrt{2}}\left[\begin{array}{c}
0 \\
v+h(\mathbf{x})
\end{array}\right] .
$$

After spontaneous symmetry breaking, we are left with only one Higgs field $h(\mathbf{x})$ in $\phi(\mathbf{x})$, after gauging away the other fields. Substituting $\phi(\mathbf{x})$ into the Lagrangian $\mathscr{L}_{3}$ (13.101), we obtain

$$
\begin{equation*}
\mathscr{L}_{3}=-\frac{G_{e}}{\sqrt{2}}\{v+h(\mathbf{x})\}\left(\bar{e}_{L} e_{R}+\bar{e}_{R} e_{L}\right) \tag{13.103}
\end{equation*}
$$

We may choose the coupling constant $G_{e}$ such that

$$
\begin{equation*}
m_{e}=\frac{G_{e} v}{\sqrt{2}} . \tag{13.104}
\end{equation*}
$$

Thus, we have generated the electron mass $m_{e}$ using the Lagrangian $\mathscr{L}_{3}$, by spontaneous symmetry breaking.

$$
\begin{equation*}
\mathscr{L}_{3}=-m_{e} \bar{e} e-\frac{m_{e}}{v} \bar{e} e h . \tag{13.105}
\end{equation*}
$$

Since the coupling constant $G_{e}$ is not known, the mass $m_{e}$ of the electron is not predicted. We have only shown that by spontaneous symmetry breaking, the electron can acquire a finite mass. The second term in Eq. (13.105) represents the interaction term of the electron with Higgs scalar. Since $v \approx 246 \mathrm{GeV}, m_{e} / v$ is negligibly small that this term will have no noticeable effect on the electro-weak interaction.

### 13.5.5 Quark masses

Just as we have generated the lepton masses, we can also generate the quark masses. But there is a difference. In the case of lepton doublet ()$_{L}$, the upper member, neutrino, has no right-handed component but only the lower member, electron, has, in addition, the right-handed component. Whereas, in the case of quark doublet ()$_{L}$, both the upper and lower members have right-handed components. This necessitates the use of two different Higgs doublets $\Phi$ and $\Phi_{c}=i \tau_{2} \Phi^{*}$ in the Lagrangian $\mathscr{L}_{4}$. The latter Higgs doublet $\Phi_{c}$ is used to generate the right-handed component for the upper member of the quark doublet.

$$
\begin{align*}
\Phi & =\left[\begin{array}{c}
\phi_{\alpha} \\
\phi_{\beta}
\end{array}\right] .  \tag{13.106}\\
\Phi_{c} & =i \tau_{2} \Phi^{*}=\left[\begin{array}{rr}
0 & 1 \\
-1 & 0
\end{array}\right]\left[\begin{array}{c}
\phi_{\alpha}^{*} \\
\phi_{\beta}^{*}
\end{array}\right]=\left[\begin{array}{c}
\phi_{\beta}^{*} \\
-\phi_{\alpha}^{*}
\end{array}\right] . \tag{13.107}
\end{align*}
$$

The Higgs doublet $\Phi_{c}$ has opposite hypercharge $(Y=-1)$ when compared to $\Phi$ with hyperchage $Y=1$, due to the special properties of $S U(2)$. So,
it is used to generate a gauge invariant Lagrangian $\mathscr{L}_{4}$ for the quarks doublet.

$$
\begin{align*}
\mathscr{L}_{4}= & -G_{d}[\bar{u}, \bar{d}]_{L}\left[\begin{array}{c}
\phi_{\alpha} \\
\phi_{\beta}
\end{array}\right] d_{R}-G_{u}[\bar{u}, \bar{d}]_{L}\left[\begin{array}{r}
\phi_{\beta}^{*} \\
-\phi_{\alpha}^{*}
\end{array}\right] u_{R} \\
& \text { +hermitian conjugate. } \tag{13.108}
\end{align*}
$$

By spontaneous symmetry breaking (ssb), $\Phi(\mathrm{x})$ and $\Phi_{c}(\mathrm{x})$ become

$$
\begin{align*}
\Phi(\mathbf{x}) & =\left[\begin{array}{c}
\phi_{\alpha} \\
\phi_{\beta}
\end{array}\right] \xrightarrow{\mathrm{ssb}} \sqrt{\frac{1}{2}}\left[\begin{array}{c}
0 \\
v+h(\mathbf{x})
\end{array}\right] ;  \tag{13.109}\\
\Phi_{c}(\mathbf{x}) & =\left[\begin{array}{c}
\phi_{\beta}^{*} \\
-\phi_{\alpha}^{*}
\end{array}\right] \xrightarrow{\text { ssb }} \sqrt{\frac{1}{2}}\left[\begin{array}{c}
v+h(\mathbf{x}) \\
0
\end{array}\right] . \tag{13.110}
\end{align*}
$$

Substituting (13.109) and (13.110) into Eq. (13.108) and simplifying, we get

$$
\begin{equation*}
\mathscr{L}_{4}=-m_{d} \bar{d} d\left(1+\frac{h}{v}\right)-m_{u} \bar{u} u\left(1+\frac{h}{v}\right), \tag{13.111}
\end{equation*}
$$

where $m_{d}=\frac{G_{d} v}{\sqrt{2}}$ and $m_{u}=\frac{G_{u} v}{\sqrt{2}}$ represent the masses of the $d$ and $u$ quarks. The quark masses depend upon the unknown coupling constants $G_{d}$ and $G_{u}$ and so they have to be treated as parameters of the theory rather than as prediction.

### 13.6 Strong interactions

## Quarks as building blocks of hadrons

Strongly interacting particles such as protons, neutrons, hyperons, pions and kaons are known as hadrons and are found to be composite objects consisting of quarks. Quarks are spin- $\frac{1}{2}$ particles carrying fractional charge and they come in six flavours. They are up, down, strange, charm, bottom and top quarks and they are presented in Table 13.2 with their quantum numbers, charges and masses.

There is a one-to-one correspondence between leptons and quarks and they are usually classified into three generations of leptons and quarks, as shown in Table 13.3 given below. There are, in total, six leptons and six flavours of quarks and the matter is built out of them.

Table 13.2: List of quarks, their flavours and their quantum numbers. Charge $Q / e=I_{3}+\frac{1}{2}(\mathrm{~B}+\mathrm{s}+\mathrm{c}+\mathrm{b}+\mathrm{t})$, where B denotes the baryon number, s the strangeness quantum number, c charm quantum number, b bottom or beauty quantum number and $t$ the top or truth quantum number.

| Flavour | B | $I$ | $I_{3}$ | s | c | b | t | $Q / e$ | Mass |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| up $(u)$ | $\frac{3}{2}$ | $\frac{1}{2}$ | $+\frac{1}{2}$ | 0 | 0 | 0 | 0 | $+\frac{2}{3}$ | $5.6 \pm 1.1 \mathrm{MeV}$ |
| down $(d)$ | $\frac{3}{2}$ | $\frac{1}{2}$ | $-\frac{1}{2}$ | 0 | 0 | 0 | 0 | $-\frac{1}{3}$ | $9.9 \pm 1.1 \mathrm{MeV}$ |
| strange $(s)$ | $\frac{3}{2}$ | 0 | 0 | -1 | 0 | 0 | 0 | $-\frac{1}{3}$ | $199 \pm 3.3 \mathrm{MeV}$ |
| charm $(c)$ | $\frac{3}{2}$ | 0 | 0 | 0 | 1 | 0 | 0 | $+\frac{2}{3}$ | $1.35 \pm 0.05 \mathrm{GeV}$ |
| bottom $(b)$ | $\frac{3}{2}$ | 0 | 0 | 0 | 0 | -1 | 0 | $-\frac{1}{3}$ | $\sim 5 \mathrm{GeV}$ |
| top $(t)$ | $\frac{3}{2}$ | 0 | 0 | 0 | 0 | 0 | 1 | $+\frac{2}{3}$ | $\sim 174 \pm 17 \mathrm{GeV}$ |

Table 13.3: Three generations of leptons and quarks

| Leptons |  |  |  | Quarks |  |  |  |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :---: |
| $\nu_{e}$ | $\nu_{\mu}$ | $\nu_{\tau}$ | $u$ | $c$ | $t$ |  |  |
| $e^{-}$ | $\mu^{-}$ | $\tau^{-}$ | $d$ | $s$ | $b$ |  |  |

## Colour degree of freedom

Soon after the proposal that quarks are the fundamental constituents of hadrons, an intense search was made to tract them down, but every attempt ended in a failure. Just as the negative result of the MichelsonMorley experiment has revolutionized the fundamental concept of physics with regard to the absolute motion and gave birth to Einstein's theory of relativity, the failure to observe the quarks in free state resulted in a new attribute of colour degree of freedom to the quarks and gave birth to the theory of quantum chromodynamics as the fundamental theory of strong interaction.

### 13.6.1 Quantum chromodynamics

The quarks are coloured objects and each quark carries three colours, say, red, green and blue and all the particles in the real world are colourless objects with a combination of these colours. It is not possible to observe the coloured quarks in free state but they certainly do exist in the bound states within a colourless particle which we observe. The evidence for the additional colour degree for quarks has come from many sources. The quarks interact through the exchange of gluons. The gluons carry colour charges and they are the vector bosons that are exchanged between coloured quarks. This is a distinguishing feature and permits gluon-gluon interaction; whereas photons which are exchanged between charged particles do not carry any charge.

Just as the charged particles interact through the exchange of photons and the quantum electrodynamics (QED) is developed as the local gauge invariance of the Lagrangian of the charged particles, the quantum chromodynamics (QCD) is also developed as the local gauge invariance of the Lagrangian of free quarks with colour degree of freedom. The appropriate gauge group is $\mathrm{SU}(3)$ which is non-abelian. The non-abelian gauge theories are generally referred to as Yang-Mills theories.

Let us now consider the free Lagrangian of the quark colour fields

$$
\begin{equation*}
\mathscr{L}_{0}=\bar{q}_{i}\left(i \gamma^{\mu} \partial_{\mu}-m\right) q_{i}, \tag{13.112}
\end{equation*}
$$

where $q_{1}, q_{2}, q_{3}$ denote the three colour fields. For simplicity, let us consider only one quark flavour. Let us consider the Lagrangian (13.112) to be invariant under local phase transformation. Since there are three colour fields, it is appropriate to consider the invariance under $\mathrm{SU}(3)$ group of phase transformation.

$$
\begin{equation*}
q(\mathbf{x})=U q(\mathbf{x})=e^{i \alpha_{a}(\mathbf{x}) T_{a}} q(\mathbf{x}) \tag{13.113}
\end{equation*}
$$

where $U$ is a $3 \times 3$ unitary matrix. $T_{a}$ with $a=1, \cdots 8$ are a set of eight independent Hermitian traceless $3 \times 3$ matrices with group parameters $\alpha_{a}$. The matrices $T_{a}$ obey the commutation relation

$$
\begin{equation*}
\left[T_{a}, T_{b}\right]_{-}=i f_{a b c} T_{c}, \tag{13.114}
\end{equation*}
$$

where $f_{a b c}$ are real constants, known as the structure constants of the group and they are antisymmetric under exchange of any pair of indices. The natural choice of matrices $T_{a}$ are $T_{a}=\lambda_{a} / 2$, where $\lambda_{a}$ denote the

Gell-Mann matrices.

$$
\begin{array}{ll}
\lambda_{1}=\left[\begin{array}{lll}
0 & 1 & 0 \\
1 & 0 & 0 \\
0 & 0 & 0
\end{array}\right], & \lambda_{2}=\left[\begin{array}{rrr}
0 & -i & 0 \\
i & 0 & 0 \\
0 & 0 & 0
\end{array}\right], \quad \lambda_{3}=\left[\begin{array}{rrr}
1 & 0 & 0 \\
0 & -1 & 0 \\
0 & 0 & 0
\end{array}\right], \\
\lambda_{4}=\left[\begin{array}{lll}
0 & 0 & 1 \\
0 & 0 & 0 \\
1 & 0 & 0
\end{array}\right], \quad \lambda_{5}=\left[\begin{array}{rrr}
0 & 0 & -i \\
0 & 0 & 0 \\
i & 0 & 0
\end{array}\right], \quad \lambda_{6}=\left[\begin{array}{lll}
0 & 0 & 0 \\
0 & 0 & 1 \\
0 & 1 & 0
\end{array}\right], \\
\lambda_{7}=\left[\begin{array}{rrr}
0 & 0 & 0 \\
0 & 0 & -i \\
0 & i & 0
\end{array}\right], \quad \lambda_{8}=\frac{1}{\sqrt{3}}\left[\begin{array}{rrr}
1 & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & -2
\end{array}\right] . \tag{13.115}
\end{array}
$$

The matrices $\lambda_{i}$ are Hermitian traceless matrices and $\operatorname{det}\left(\lambda_{i}\right)=1$. They obey the commutation relations

$$
\left[\lambda_{a}, \lambda_{b}\right]_{-}=2 i f_{a b c} \lambda_{c}
$$

where $f_{a b c}$ are structure constants given by

$$
\begin{align*}
& f_{123}=1, \\
& f_{458}=f_{678}=\frac{\sqrt{3}}{2}  \tag{13.116}\\
& f_{147}=f_{165}=f_{246}=f_{257}=f_{345}=f_{376}=\frac{1}{2}
\end{align*}
$$

They are antisymmetric under interchange of any pair of indices. The structure constants with other indices are zero.

Let us consider infinitesimal phase transformations and find the conditions under which the Lagrangian $\mathscr{L}$ is invariant under $\mathrm{SU}(3)$ local gauge transformation.

$$
\begin{align*}
q(\mathbf{x}) & \rightarrow\left\{1+i \alpha_{a}(\mathbf{x}) T_{a}\right\} q(\mathbf{x})  \tag{13.117}\\
\partial_{\mu} q & \rightarrow\left\{1+i \alpha_{a} T_{a}\right\} \partial_{\mu} q+i T_{a} q \partial_{\mu} \alpha_{a} \tag{13.118}
\end{align*}
$$

The extra term in Eq. (13.118) spoils the gauge invariance of $\mathscr{L}$ and let us proceed to restore the gauge invariance by introducing eight gauge fields $G_{\mu}^{a}, a=1, \cdots, 8$, each transforming under gauge transformation as shown below:

$$
\begin{equation*}
G_{\mu}^{a} \rightarrow G_{\mu}^{a}-\frac{1}{g} \partial_{\mu} \alpha_{a} \tag{13.119}
\end{equation*}
$$

This is the procedure, we adopted in QED and let us follow it here and define a covariant derivative

$$
\begin{equation*}
D_{\mu}=\partial_{\mu}+i g T_{a} G_{\mu}^{a} \tag{13.120}
\end{equation*}
$$

and use it to replace $\partial_{\mu}$ by $D_{\mu}$ in the Lagrangian $\mathscr{L}_{0}$, given by Eq. (13.112).

$$
\begin{equation*}
\mathscr{L}=\bar{q}\left(i \gamma_{\mu} D_{\mu}-m\right) q=\bar{q}\left(i \gamma_{\mu} \partial_{\mu}-m\right) q-g\left(\bar{q} \gamma^{\mu} T_{a} q\right) G_{\mu}^{a} . \tag{13.121}
\end{equation*}
$$

The Lagrangian (13.121) is analogous to the QED Lagrangian but it is not gauge-invariant in the case of QCD due to the non-abelian nature of the gauge transformation. Under infinitesimal $\mathrm{SU}(3)$ gauge transformation,

$$
\begin{align*}
\bar{q} \gamma^{\mu} T_{a} q \rightarrow \bar{q}^{\prime} \gamma^{\mu} T_{a} q^{\prime} & =\left(1-i \alpha_{b} T_{b}\right) \bar{q} \gamma^{\mu} T_{a} q\left(1+i \alpha_{b} T_{b}\right) \\
& =\bar{q} \gamma^{\mu} T_{a} q+i \alpha_{b} \bar{q} \gamma^{\mu}\left[T_{a}, T_{b}\right]_{-} q \\
& =\bar{q} \gamma^{\mu} T_{a} q-f_{a b c} \alpha_{b}\left(\bar{q} \gamma^{\mu} T_{c} q\right) . \tag{13.122}
\end{align*}
$$

The result (13.122) is obtained using the commutation relation (13.114) and it suggests that the Lagrangian (13.121) can be made invariant by modifying the gauge transformation property (13.119) of the gauge boson.

$$
\begin{equation*}
G_{\mu}^{a} \rightarrow G_{\mu}^{a}-\frac{1}{g} \partial_{\mu} \alpha_{a}-f_{a b c} \alpha_{b} G_{\mu}^{c} \tag{13.123}
\end{equation*}
$$

By adding the kinetic energy term for each of the gauge fields $G_{\mu}^{a}$ to the Lagrangian (13.121), we obtain the total gauge invariant QCD Lagrangian.

$$
\begin{equation*}
\mathscr{L}=\bar{q}\left(i \gamma_{\mu} \partial_{\mu}-m\right) q-g\left(\bar{q} \gamma^{\mu} T_{a} q\right) G_{\mu}^{a}-\frac{1}{4} G_{\mu \nu}^{a} G_{a}^{\mu \nu} \tag{13.124}
\end{equation*}
$$

The additional term involves $G_{\mu \nu}^{a}$ which has a more complicated structure than its counterpart in QED.

$$
\begin{equation*}
G_{\mu \nu}^{a}=\partial_{\mu} G_{\nu}^{a}-\partial_{\nu} G_{\mu}^{a}-g f_{a b c} G_{\mu}^{b} G_{\nu}^{c} \tag{13.125}
\end{equation*}
$$

The QCD Lagrangian (13.124) gives the interaction between the coloured quarks with the vector bosons (gluons) and also between the vector bosons themselves. The vector bosons are massless just as photons but they carry colour charges unlike photons.

The Lagrangian (13.124) involves the following terms

$$
\bar{q} q, \bar{q} q G, G^{2}, G^{3} \text { and } G^{4}
$$

indicating the propagation of quark, interaction of quark with gluon, gluon propagation, three gluon vertex and four gluon vertex, which can be represented diagrammatically as shown in Fig. 13.3 below:

Diagrams (d) and (e) represent the cubic and quartic interactions between gluons because gluons carry colour charges. Such interactions are absent in the case of e.m. interactions because photons do not carry electric charges.


Figure 13.3: Diagrammatic representation of the various terms in the QCD Lagrangian (13.124).

### 13.7 Summary

Let us now summarize the results of our study. The requirement of invariance of Lagrangian under a local gauge transformation leads to the study of fundamental interactions of particles. Take for example, the Lagrangian for a free electron. The requirement that the Lagrangian should be locally invariant under the same type of transformation can only be fulfilled by introducing additional fields, known as gauge fields. In a sense, the electromagnetic field is a consequence of the local gauge symmetry of the Lagrangian for the electron.

This procedure can be extended to the study of more complex gauge transformations to unravel other types of interactions such as weak interactions and strong interactions. The weak interactions are mediated by intermediate vector bosons which are massive but the gauge-invariant Lagrangian does not allow mass term for the gauge bosons. This difficulty is overcome by introducing the concept of spontaneous symmetry breaking and the usage of Higgs mechanism. The choice of negative value for $\mu^{2}$ in the potential term of the Lagrangian shifts the ground state from the origin to a finite value of $\phi$. Thus the phase (gauge) symmetry of the Lagrangian is not shared by the ground state (vacuum state). This phenomenon is known as spontaneous symmetry breaking.

The breaking of the phase (gauge) symmetry by the ground state enables the gauge boson to acquire mass. At the same time, the fact that this symmetry is retained by the Lagrangian enables the renormalization to go through. This last feature is found to be very important because it ensures a systematic cancellation of infinities, without which there will not be any predictable field theory. The renormalizability of electro-weak gauge theory in the presence of spontaneous symmetry breaking was shown to be
valid by $\mathrm{t}^{\prime} \mathrm{Hooft}^{8}$ for which he was awarded the Nobel Prize in 1999.
Thus, the gauge principle provides a general scheme for introducing interaction by constructing gauge field theories. To this end one starts with a Lagrangian for a matter field and derives the interaction by introducing exactly those fields that make the Lagrangian invariant under a relevant local gauge transformation. It seems that all fundamental forces can be described by such local gauge field theories. Gauge symmetry plays a crucial role in determining the dynamics of the theory since the nature of gauge transformation determines the possible interaction. The structure of these transformations are characterized by special mathematical groups: $\mathrm{U}(1)$ for $\mathrm{QED}, \mathrm{SU}(2) \times \mathrm{U}(1)$ for electroweak interaction, $\mathrm{SU}(3)$ for strong interaction. The relations between these groups are exploited in programs for the unification of the fundamental types of interaction. From a more technical point of view, gauge symmetries are important tools in proofs of renormalizability. So, it turns out that gauge invariance plays as much an important role as Lorentz invariance in the selection of theories. Since gauge invariance plays a pivotal role in the discovery of quantum field theories, it is a paradigm case for how a rich mathematical structure can help in the construction of theories.

The gauge theories invoke the scalar fields, known as the Higgs fields, and are based on the concepts of spontaneous symmetry breaking and the Higgs mechanism for generating masses for the gauge bosons, leptons and quarks. The outcome of these theories is the additional production of Higgs boson of mass $m_{h}=\sqrt{2 \lambda v^{2}}$, which can be considered as the quantum excitation of of one of the four components of Higgs field. Hence its discovery is considered to be very important for verifying the validity of gauge theories. That is why the CERN undertook the arduous task of constructing the most expensive and complex experimental facility, the Large Hadron Collider (LHC) and inviting international collaboration for conducting the experiments with LHC; which finally resulted in the discovery of Higgs boson with a mass of about 126 GeV . The discovery of Higgs boson was confirmed by CERN on March 14, 2013, following which Peter Higgs and Francois Englert were jointly awarded the Nobel Prize in Physics for the year 2013.

[^79]
## Review Questions

13.1 (a) State and explain Noether's theorem
(b) Show that the global gauge invariance of Dirac's Lagrangian leads to conservation of current.
13.2 Show that the Dirac Lagrangian density is invariant under global phase transformation but not under local phase transformation. In order to make it invariant under local phase transformation, one has to introduce a gauge field $A^{\mu}$ which represents the e.m. field.
13.3 Write down the Lagrangian for a real scalar field $\phi(\mathbf{x})$ and show that it has reflection symmetry. Explain the concept of spontaneous symmetry breaking by introducing a new field $\eta(\mathbf{x})$, in terms of which the Lagrangian has lost the reflection symmetry. Find the mass of the new scalar field $\eta(\mathbf{x})$.
13.4 Write down the Lagrangian of a complex scalar field $\phi(\mathbf{x})$ which is invariant under global $U(1)$ gauge transformation. By spontaneous symmetry breaking, show that one can obtain a massive scalar field known as gauge boson and a massless scalar field known as Goldstone boson.
13.5 Consider the Lagrangian of a charged scalar field of mass $\mu$, which is invariant under a $U(1)$ local gauge transformation. Show that, by spontaneous symmetry breaking, the gauge boson acquires a mass and a massless Goldstone boson is also generated. Show how, by using the unitary gauge, the massless Goldstone boson is eliminated.
13.6 Discuss Higg's mechanism of giving mass to the gauge boson by absorbing the massless Goldstone boson.
13.7 Consider the Lagrangian of a $S U(2)$ doublet complex scalar field which is invariant under global $S U(2)$ phase transformation. Discuss how you can make it invariant under local gauge transformation by introducing three vector gauge fields and how, by invoking Higg's mechanism, you can make them acquire masses.
13.8 Discuss the $\mathrm{SU}(2) \times \mathrm{U}(1)$ group of gauge theory of electro-weak interactions as developed by Weinberg and Salam and explain how the intermediate vector bosons that mediate the weak interactions acquire masses whereas the photon that mediates the e.m. interaction remains massless.
13.9 Explain why in electro-weak theory, the mass terms of the leptons in the Lagrangian have to be omitted to make it gauge-invariant and how by invoking Higgs fields, the lepton masses can be generated by Higgs mechanism.
13.10 Discuss why in electro-weak theory, the quarks have to be massless for the Lagrangian to be gauge-invariant and how the quarks can be made to acquire masses by spontaneous symmetry breaking.
13.11 Discuss briefly how the quantum chromodynamics can be treated as the theory of strong interactions.
13.12 Discuss the $\operatorname{SU}(3)$ gauge theory of QCD and obtain the appropriate Lagrangian density which is gauge-invariant. Obtain expressions for the gauge field and also for the free field Lagrangian of the gauge boson.
13.13 Discuss how the gauge theory is considered as the standard model of elementary particles and how does the discovery of Higgs boson vindicate the gauge theory of elementary particles.

## Problems

13.1 Consider the following Lagrangian density of a complex scalar field:

$$
\mathscr{L}=\left(\partial_{\mu} \phi\right)^{*}\left(\partial^{\mu} \phi\right)-m^{2} \phi^{*} \phi .
$$

Show that the $U(1)$ phase invariance of the Lagrangian implies the existence of a conserved current.

$$
j^{\mu}=-i e\left(\phi^{*} \partial^{\mu} \phi-\phi \partial^{\mu} \phi^{*}\right) .
$$

13.2 Show that $U(\boldsymbol{\alpha})=e^{i \alpha_{i} \tau_{i} / 2}$, where $\alpha_{i}$ represents a real number and $\tau_{i}$ the Pauli matrices, is a unitary matrix. Write down explicitly the operator $U(\boldsymbol{\alpha})$ as a $2 \times 2$ unitary matrix. Show that $\operatorname{det}(U(\boldsymbol{\alpha}))=1$ and hence $U(\boldsymbol{\alpha})$ represents a $\mathrm{SU}(2)$ transformation.
13.3 Show that the Lagrangian of a $\operatorname{SU}(2)$ doublet complex scalar field $\phi$ is invariant under an infinitesimal local gauge transformation with $\boldsymbol{\alpha}(\mathbf{x})$,

$$
\phi \rightarrow \phi^{\prime}=e^{i \alpha_{a} \tau_{a} / 2} \phi,
$$

if the derivative $\partial_{\mu}$ is replaced by the covariant derivative

$$
D_{\mu}=\partial_{\mu}+\frac{1}{2} i g \tau_{a} W_{\mu}^{a},
$$

by invoking three gauge vector fields $W_{\mu}^{a}$ with $a=1,2,3$, which transform as

$$
\boldsymbol{W}_{\mu} \rightarrow \boldsymbol{W}_{\mu}^{\prime}=\boldsymbol{W}_{\mu}-\frac{1}{g} \partial_{\mu} \boldsymbol{\alpha}-\boldsymbol{\alpha} \times \boldsymbol{W}_{\mu} .
$$

The strength of the $\mathrm{SU}(2)$ coupling to the gauge field is denoted by $g$ and $\tau_{a}$ with $a=1,2,3$ denote the three iso-spin Pauli matrices.
13.4 Consider the rotation of a vector $\boldsymbol{A}$ through an infinitesimal angle $\boldsymbol{\alpha}$ in a three-dimensional space and show that the rotated vector $\boldsymbol{A}^{\prime}$ can be expressed as $\boldsymbol{A}^{\prime}=\boldsymbol{A}-\boldsymbol{\alpha} \times \boldsymbol{A}$ such that the scalar product $\boldsymbol{A}^{\prime} \cdot \boldsymbol{A}^{\prime}=\boldsymbol{A} \cdot \boldsymbol{A}$ is invariant under rotation.
13.5 Show that under $\operatorname{SU}(2)$ non-abelian gauge transformation, the quantity

$$
\mathcal{W}_{\mu \nu}=\partial_{\mu} \boldsymbol{W}_{\nu}-\partial_{\nu} \boldsymbol{W}_{\mu}
$$

does not transform as an iso-vector but the quantity

$$
\boldsymbol{W}_{\mu \nu}=\mathcal{W}_{\mu \nu}-g\left(\boldsymbol{W}_{\mu} \times \boldsymbol{W}_{\nu}\right)
$$

transforms as an iso-vector.
13.6 Show that the free Lagrangian density of the produced gauge fields

$$
\mathscr{L}_{g}=-\frac{1}{4} \boldsymbol{W}_{\mu \nu}^{a} \boldsymbol{W}_{a}^{\mu \nu}, \quad(a=1,2,3)
$$

is invariant under $\mathrm{SU}(2)$ non-abelian gauge transformation.
13.7 From Fermi's theory of muon decay, the Fermi coupling constant $G_{F}$ is found to be

$$
G_{F}=1.166 \times 10^{-5} \mathrm{Gev}^{-2} .
$$

Comparing Fermi's theory with Weinberg-Salam Model, we obtain a relation

$$
\frac{G_{F}}{\sqrt{2}}=\frac{g^{2}}{8 M_{W}^{2}}=\frac{1}{2 v^{2}} .
$$

From the $\nu_{\mu}-e$ scattering cross section

$$
\nu_{\mu}+e^{-} \rightarrow \nu_{\mu}+e^{-}
$$

one gets a value for the mixing angle (Weinberg angle) $\theta_{W}$.

$$
\sin ^{2} \theta_{W}=0.225
$$

Given the fine structure constant $\frac{e^{2}}{4 \pi}=\frac{1}{137}$, find the values of $v$, the fundamental coupling constants $g$ and $g^{\prime}$ and the masses of the gauge bosons $M_{W}$ and $M_{Z}$.

## Solutions to Problems

13.1 Under $U(1)$ phase invariance,

$$
\phi(\mathbf{x}) \rightarrow e^{i \alpha} \phi(\mathbf{x}) ; \quad \phi^{*}(\mathbf{x}) \rightarrow e^{-i \alpha} \phi^{*}(\mathbf{x}) .
$$

Under infinitesimal $U(1)$ phase invariance,

$$
\phi(\mathbf{x}) \rightarrow(1+i \alpha) \phi(\mathbf{x}) ; \quad \phi^{*}(\mathbf{x}) \rightarrow(1-i \alpha) \phi^{*}(\mathbf{x}) .
$$

Thus, we get

$$
\begin{aligned}
& \delta \phi(\mathbf{x})=i \alpha \phi(\mathbf{x}) ; \\
& \delta \phi^{*}(\mathbf{x})=-i \alpha \phi^{*}(\mathbf{x}) . \\
& \delta\left(\partial_{\mu} \phi(\mathbf{x})\right)=i \alpha\left(\partial_{\mu} \phi(\mathbf{x})\right) ; \delta\left(\partial_{\mu} \phi^{*}(\mathbf{x})\right)=-i \alpha\left(\partial_{\mu} \phi^{*}(\mathbf{x})\right) .
\end{aligned}
$$

Invariance of $\mathscr{L}$ implies $\delta \mathscr{L}=0$ under the $U(1)$ phase transformation. Following the method given in Sec. 13.1 to get Eq. (13.12), we obtain

$$
\delta \mathscr{L}=i \alpha \partial_{\mu}\left[\frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \phi\right)} \phi-\phi^{*} \frac{\partial \mathscr{L}}{\partial\left(\partial_{\mu} \phi^{*}\right)}\right]=0 .
$$

Using the given Lagrangian,

$$
\mathscr{L}=\left(\partial_{\mu} \phi\right)^{*}\left(\partial^{\mu} \phi\right)-m^{2} \phi^{*} \phi
$$

we obtain

$$
\delta \mathscr{L}=i \alpha \partial_{\mu}\left[\left(\partial^{\mu} \phi\right)^{*} \phi-\phi^{*}\left(\partial^{\mu} \phi\right)\right]=\partial_{\mu} j^{\mu}=0 .
$$

By replacing $\alpha$ by $e$, we obtain the conserved current.

$$
j^{\mu}=-i e\left(\phi^{*} \partial^{\mu} \phi-\phi\left(\partial^{\mu} \phi\right)^{*}\right)
$$

13.2 First let us show that $U(\boldsymbol{\alpha})$ is a unitary matrix. $U(\boldsymbol{\alpha})^{\dagger}=U(\boldsymbol{\alpha})^{-1}$.

Since $\alpha_{i}$ is real and $\tau_{i}$ is a Pauli matrix which is Hermitian $\left(\tau_{i}^{\dagger}=\tau_{i}\right)$, we obtain

$$
U(\boldsymbol{\alpha})^{\dagger}=e^{-i \alpha_{i} \tau_{i}^{\dagger} / 2}=e^{-i \alpha_{i} \tau_{i} / 2}=U(\boldsymbol{\alpha})^{-1}
$$

Let us now show explicitly that $U_{x}$ is a $2 \times 2$ matrix.

$$
\begin{aligned}
U_{x} & =e^{i \alpha_{x} \tau_{x} / 2} \\
& =1+i \frac{\alpha_{x}}{2} \tau_{x}-\frac{\left\{\left(\alpha_{x} / 2\right) \tau_{x}\right\}^{2}}{2!}-i \frac{\left\{\left(\alpha_{x} / 2\right) \tau_{x}\right\}^{3}}{3!}+\frac{\left\{\left(\alpha_{x} / 2\right) \tau_{x}\right\}^{4}}{4!}+\cdots \\
& =1+i\left(\alpha_{x} / 2\right) \tau_{x}-\frac{\left(\alpha_{x} / 2\right)^{2}}{2!}-i \frac{\left(\alpha_{x} / 2\right)^{3}}{3!} \tau_{x}+\frac{\left(\alpha_{x} / 2\right)^{4}}{4!}+\cdots \\
& =\left(1-\frac{\left(\alpha_{x} / 2\right)}{2!}+\frac{\left(\alpha_{x} / 2\right)^{4}}{4!}-\cdots\right)+i \tau_{x}\left(\alpha_{x} / 2-\frac{\left(\alpha_{x} / 2\right)^{3}}{3!}+\cdots\right) \\
& =\cos \frac{\alpha_{x}}{2}+i \tau_{x} \sin \frac{\alpha_{x}}{2} \\
& =\left[\begin{array}{cc}
\cos \frac{\alpha_{x}}{2} & i \sin \frac{\alpha_{x}}{2} \\
i \sin \frac{\alpha_{x}}{2} & \cos \frac{\alpha_{x}}{2}
\end{array}\right]
\end{aligned}
$$

Above, we have used the property that $\tau_{x}^{2}=1$. It can be easily seen that $\operatorname{det} U_{x}=1$. By the same procedure, we can obtain the other components $U_{y}$ and $U_{z}$. Since $U=U_{x} U_{y} U_{z}$ and $\operatorname{det} U=\operatorname{det} U_{x} \operatorname{det} U_{y} \operatorname{det} U_{z}$, we obtain $\operatorname{det} U=1$. So, $U(\boldsymbol{\alpha})$ is a unitary $2 \times 2$ matrix with $\operatorname{det} U=1$. So, it represents a $\mathrm{SU}(2)$ unitary transformation.
13.3 Under infinitesimal gauge transformation,

$$
\psi \rightarrow \psi^{\prime}=e^{(i / 2) \boldsymbol{\tau} \cdot \boldsymbol{\alpha}(\mathbf{x})} \psi=\left(1+\frac{i}{2} \boldsymbol{\tau} \cdot \boldsymbol{\alpha}\right) \psi
$$

The covariant derivative transforms in the same way as the field $\psi$.

$$
\begin{equation*}
D \psi \rightarrow(D \psi)^{\prime}=\left(1+\frac{i}{2} \boldsymbol{\tau} \cdot \boldsymbol{\alpha}\right)\left(\partial_{\mu}+\frac{i g}{2} \boldsymbol{\tau} \cdot \boldsymbol{W}\right) \psi \tag{13.126}
\end{equation*}
$$

Under infinitesimal gauge transformation, if the gauge field $\boldsymbol{W}$ transforms as

$$
\boldsymbol{W} \rightarrow \boldsymbol{W}^{\prime}=\boldsymbol{W}+\delta \boldsymbol{W}
$$

then

$$
\begin{equation*}
D \psi \rightarrow D^{\prime} \psi^{\prime}=\left\{\partial_{\mu}+\frac{i g}{2} \boldsymbol{\tau} \cdot(\boldsymbol{W}+\delta \boldsymbol{W})\right\}\left(1+\frac{i}{2} \boldsymbol{\tau} \cdot \boldsymbol{\alpha}\right) \psi \tag{13.127}
\end{equation*}
$$

Equating the right hand sides of Eqs. (13.127) and (13.126) and retaining only the first order terms, we get

$$
\begin{aligned}
\frac{i}{2} \partial_{\mu}(\boldsymbol{\tau} \cdot \boldsymbol{\alpha}) \psi+ & \frac{i g}{2} \boldsymbol{\tau} \cdot(\boldsymbol{W}+\delta \boldsymbol{W}) \psi-\frac{g}{4}(\boldsymbol{\tau} \cdot \boldsymbol{W})(\boldsymbol{\tau} \cdot \boldsymbol{\alpha}) \psi \\
& =\frac{i g}{2} \boldsymbol{\tau} \cdot \boldsymbol{W} \psi-\frac{g}{4}(\boldsymbol{\tau} \cdot \boldsymbol{\alpha})(\boldsymbol{\tau} \cdot \boldsymbol{W}) \psi
\end{aligned}
$$

The above equation simplifies to

$$
\begin{align*}
(\boldsymbol{\tau} \cdot \delta \boldsymbol{W}) \psi & =-\frac{1}{g} \partial_{\mu}(\boldsymbol{\tau} \cdot \boldsymbol{\alpha}) \psi+\frac{i}{2}\{(\boldsymbol{\tau} \cdot \boldsymbol{\alpha})(\boldsymbol{\tau} \cdot \boldsymbol{W})-(\boldsymbol{\tau} \cdot \boldsymbol{W})(\boldsymbol{\tau} \cdot \boldsymbol{\alpha})\} \psi \\
& =-\frac{1}{g} \partial_{\mu}(\boldsymbol{\tau} \cdot \boldsymbol{\alpha}) \psi+\frac{i}{2}[\boldsymbol{\tau} \cdot \boldsymbol{\alpha}, \boldsymbol{\tau} \cdot \boldsymbol{W}]_{-} \psi \tag{13.128}
\end{align*}
$$

Let us now evaluate the commutator $[\boldsymbol{\tau} \cdot \boldsymbol{\alpha}, \boldsymbol{\tau} \cdot \boldsymbol{W}]_{-}$, using the familiar relation involving Pauli matrices.

$$
(\boldsymbol{\sigma} \cdot \boldsymbol{A})(\boldsymbol{\sigma} \cdot \boldsymbol{B})=\boldsymbol{A} \cdot \boldsymbol{B}+i \boldsymbol{\sigma} \cdot(\boldsymbol{A} \times \boldsymbol{B})
$$

It follows that

$$
\begin{aligned}
(\boldsymbol{\tau} \cdot \boldsymbol{\alpha})(\boldsymbol{\tau} \cdot \boldsymbol{W}) & =\boldsymbol{\alpha} \cdot \boldsymbol{W}+i \boldsymbol{\tau} \cdot(\boldsymbol{\alpha} \times \boldsymbol{W}) \\
(\boldsymbol{\tau} \cdot \boldsymbol{W})(\boldsymbol{\tau} \cdot \boldsymbol{\alpha}) & =\boldsymbol{W} \cdot \boldsymbol{\alpha}+i \boldsymbol{\tau} \cdot(\boldsymbol{W} \times \boldsymbol{\alpha}) \\
& =\boldsymbol{\alpha} \cdot \boldsymbol{W}-i \boldsymbol{\tau} \cdot(\boldsymbol{\alpha} \times \boldsymbol{W})
\end{aligned}
$$

Substituting the value of the commutator

$$
\begin{equation*}
[\boldsymbol{\tau} \cdot \boldsymbol{\alpha}, \boldsymbol{\tau} \cdot \boldsymbol{W}]_{-}=2 i \boldsymbol{\tau} \cdot(\boldsymbol{\alpha} \times \boldsymbol{W}) \tag{13.129}
\end{equation*}
$$

in Eq. (13.128), we obtain

$$
(\boldsymbol{\tau} \cdot \delta \boldsymbol{W}) \psi=-\frac{1}{g} \partial_{\mu}(\boldsymbol{\tau} \cdot \boldsymbol{\alpha}) \psi-\boldsymbol{\tau} \cdot(\boldsymbol{\alpha} \times \boldsymbol{W}) \psi
$$

This yields the transformation formula for the gauge field $\boldsymbol{W}_{\mu}$ under $\mathrm{SU}(2)$ gauge transformation.

$$
\boldsymbol{W}_{\mu} \rightarrow \boldsymbol{W}_{\mu}^{\prime}=\boldsymbol{W}_{\mu}+\delta \boldsymbol{W}_{\mu}=\boldsymbol{W}_{\mu}-\frac{1}{g} \partial_{\mu} \boldsymbol{\alpha}-\boldsymbol{\alpha} \times \boldsymbol{W}_{\mu}
$$

## Alternative Method

Under infinitesimal gauge transformation $U=1+(i / 2) \boldsymbol{\tau} \cdot \boldsymbol{\alpha}$, the field $\psi$ and its covariant derivative $D \psi$ transform as given below:

$$
\psi \rightarrow \psi^{\prime}=U \psi ; \quad D \psi \rightarrow D^{\prime} \psi^{\prime}=U D U^{\dagger} U \psi
$$

Under infinitesimal gauge transformation, the covariant derivative $D \psi=$ $\partial_{\mu}+(i g / 2) \boldsymbol{\tau} \cdot \boldsymbol{W}$ becomes (retaining only the first order terms in $\boldsymbol{\alpha}$ since we are considering an infinitesimal transformation)

$$
\begin{align*}
D^{\prime}= & U D U^{\dagger}=\left(1+\frac{i}{2} \boldsymbol{\tau} \cdot \boldsymbol{\alpha}\right)\left(\partial_{\mu}+\frac{i g}{2} \boldsymbol{\tau} \cdot \boldsymbol{W}\right)\left(1-\frac{i}{2} \boldsymbol{\tau} \cdot \boldsymbol{\alpha}\right) \\
= & \left(1+\frac{i}{2} \boldsymbol{\tau} \cdot \boldsymbol{\alpha}\right)\left(\partial_{\mu}+\frac{i g}{2} \boldsymbol{\tau} \cdot \boldsymbol{W}\right)-\frac{i}{2}\left\{\boldsymbol{\tau} \cdot\left(\partial_{\mu} \boldsymbol{\alpha}\right)+(\boldsymbol{\tau} \cdot \boldsymbol{\alpha}) \partial_{\mu}\right\} \\
& +\frac{g}{4}(\boldsymbol{\tau} \cdot \boldsymbol{W})(\boldsymbol{\tau} \cdot \boldsymbol{\alpha}) \\
= & \partial_{\mu}+\frac{i g}{2} \boldsymbol{\tau} \cdot \boldsymbol{W}-\frac{i}{2} \boldsymbol{\tau} \cdot\left(\partial_{\mu} \boldsymbol{\alpha}\right)-\frac{g}{4}[(\boldsymbol{\tau} \cdot \boldsymbol{\alpha}),(\boldsymbol{\tau} \cdot \boldsymbol{W})]_{-} \tag{13.130}
\end{align*}
$$

Substituting the value of the commutator given by Eq. (13.129), we get

$$
\begin{align*}
D^{\prime} & =\partial_{\mu}+\frac{i g}{2} \boldsymbol{\tau} \cdot\left\{\boldsymbol{W}-\frac{1}{g} \partial_{\mu} \boldsymbol{\alpha}-\boldsymbol{\alpha} \times \boldsymbol{W}\right\} \\
& =\partial_{\mu}+\frac{i g}{2} \boldsymbol{\tau} \cdot \boldsymbol{W}^{\prime} \tag{13.131}
\end{align*}
$$

Thus we obtain the transformation law for the gauge boson $\boldsymbol{W}$.

$$
\boldsymbol{W}^{\prime}=\boldsymbol{W}-\frac{1}{g} \partial_{\mu} \boldsymbol{\alpha}-\boldsymbol{\alpha} \times \boldsymbol{W}
$$

13.4 Consider the rotation of a vector $\boldsymbol{A}$ in the x-y plane through an infinitesimal angle $\alpha_{z}$ about the z - axis, as shown in Fig. The rotated vector is denoted by $\boldsymbol{A}^{\prime}$.

$$
\begin{aligned}
& A_{x}=A \cos \theta, \quad A_{y}=A \sin \theta, \quad A_{z}=0 \\
A_{x}^{\prime}= & A \cos \left(\theta-\alpha_{z}\right)=A\left(\cos \theta \cos \alpha_{z}+\sin \theta \sin \alpha_{z}\right)=A\left(\cos \theta+\sin \theta \alpha_{z}\right) \\
= & A_{x}+A_{y} \alpha_{z} . \\
A_{y}^{\prime}= & A \sin \left(\theta-\alpha_{z}\right)=A\left(\sin \theta \cos \alpha_{z}-\cos \theta \sin \alpha_{z}\right)=A\left(\sin \theta-\cos \theta \alpha_{z}\right) \\
= & A_{y}-A_{x} \alpha_{z} . \\
A_{z}^{\prime}= & A_{z} .
\end{aligned}
$$



Figure 13.4: The rotation of vector $\boldsymbol{A}$ by an infinitesimal angle $\alpha_{z}$ about the z-axis. The rotated vector is $\boldsymbol{A}^{\prime}$.

In a similar way, we can consider infinitesimal rotations of the vector $\boldsymbol{A}$ through an angle $\alpha_{x}$ about the $\mathrm{x}-$ axis and an angle $\alpha_{y}$ about the y - axis and obtain the components of the rotated vector, $\boldsymbol{A}^{\prime}$.

$$
\begin{array}{ll}
A_{x}^{\prime}=A_{x}, & A_{y}^{\prime}=A_{y}+\alpha_{x} A_{z},
\end{array} \quad A_{z}^{\prime}=A_{z}-\alpha_{x} A_{y} .
$$

Collecting the above results, we can express rotated vector in a compact form:

$$
\begin{equation*}
A^{\prime}=A-\alpha \times A \tag{13.132}
\end{equation*}
$$

Retaining only the first order terms in $\alpha$, it can be easily seen that

$$
\boldsymbol{A}^{\prime} \cdot \boldsymbol{A}^{\prime}=(\boldsymbol{A}-\boldsymbol{\alpha} \times \boldsymbol{A}) \cdot(\boldsymbol{A}-\boldsymbol{\alpha} \times \boldsymbol{A})=\boldsymbol{A} \cdot \boldsymbol{A}
$$

since $(\boldsymbol{\alpha} \times \boldsymbol{A}) \cdot \boldsymbol{A}=\boldsymbol{A} \cdot(\boldsymbol{\alpha} \times \boldsymbol{A})=\boldsymbol{\alpha} \cdot(\boldsymbol{A} \times \boldsymbol{A})=0$.
13.5 Under $\mathrm{SU}(2)$ gauge transformation,

$$
\begin{align*}
\mathcal{W}_{\mu \nu} \rightarrow \mathcal{W}_{\mu \nu}^{\prime}= & \partial_{\mu} \boldsymbol{W}_{\nu}^{\prime}-\partial_{\nu} \boldsymbol{W}_{\mu}^{\prime} \\
= & \partial_{\mu}\left(\boldsymbol{W}_{\nu}-\frac{1}{g} \partial_{\nu} \boldsymbol{\alpha}-\boldsymbol{\alpha} \times \boldsymbol{W}_{\nu}\right) \\
& -\partial_{\nu}\left(\boldsymbol{W}_{\mu}-\frac{1}{g} \partial_{\mu} \boldsymbol{\alpha}-\boldsymbol{\alpha} \times \boldsymbol{W}_{\mu}\right) \\
= & \partial_{\mu} \boldsymbol{W}_{\nu}-\partial_{\nu} \boldsymbol{W}_{\mu}-\partial_{\mu}\left(\boldsymbol{\alpha} \times \boldsymbol{W}_{\nu}\right)+\partial_{\nu}\left(\boldsymbol{\alpha} \times \boldsymbol{W}_{\mu}\right) \\
= & \partial_{\mu} \boldsymbol{W}_{\nu}-\partial_{\nu} \boldsymbol{W}_{\mu}-\boldsymbol{\alpha} \times\left(\partial_{\mu} \boldsymbol{W}_{\nu}-\partial_{\nu} \boldsymbol{W}_{\mu}\right) \\
& -\left(\partial_{\mu} \boldsymbol{\alpha} \times \boldsymbol{W}_{\nu}-\partial_{\nu} \boldsymbol{\alpha} \times \boldsymbol{W}_{\mu}\right) \tag{13.133}
\end{align*}
$$

In Eq. (13.133), if the last term was absent, then the quantity $\mathcal{W}_{\mu \nu}$ can be regarded as a vector in iso-spin space as per the results of Prob. 13.4. So, let us try to eliminate the last term by adding or subtracting an additional term.
Let us consider the quantity $g\left(\boldsymbol{W}_{\mu} \times \boldsymbol{W}_{\nu}\right)$. Under $\mathrm{SU}(2)$ gauge transformation, it transforms as

$$
\begin{align*}
g\left(\boldsymbol{W}_{\mu} \times \boldsymbol{W}_{\nu}\right) \rightarrow & g\left(\boldsymbol{W}_{\mu}^{\prime} \times \boldsymbol{W}_{\nu}^{\prime}\right) \\
g\left(\boldsymbol{W}_{\mu}^{\prime} \times \boldsymbol{W}_{\nu}^{\prime}\right)= & g\left(\boldsymbol{W}_{\mu}-\frac{1}{g} \partial_{\mu} \boldsymbol{\alpha}-\boldsymbol{\alpha} \times \boldsymbol{W}_{\mu}\right) \\
& \times\left(\boldsymbol{W}_{\nu}-\frac{1}{g} \partial_{\nu} \boldsymbol{\alpha}-\boldsymbol{\alpha} \times \boldsymbol{W}_{\nu}\right) \\
= & g\left(\left(\boldsymbol{W}_{\mu} \times \boldsymbol{W}_{\nu}\right)-\boldsymbol{W}_{\mu} \times \partial_{\nu} \boldsymbol{\alpha}-\partial_{\mu} \boldsymbol{\alpha} \times \boldsymbol{W}_{\nu}\right. \\
& -g\left\{\boldsymbol{W}_{\mu} \times\left(\boldsymbol{\alpha} \times \boldsymbol{W}_{\nu}\right)+\left(\boldsymbol{\alpha} \times \boldsymbol{W}_{\mu}\right) \times \boldsymbol{W}_{\nu}\right\} \tag{13.134}
\end{align*}
$$

Let us consider the vector identity

$$
(\boldsymbol{A} \times \boldsymbol{B}) \times \boldsymbol{C}+(\boldsymbol{B} \times \boldsymbol{C}) \times \boldsymbol{A}+(\boldsymbol{C} \times \boldsymbol{A}) \times \boldsymbol{B}=0
$$

which on rearrangement yields

$$
(\boldsymbol{A} \times \boldsymbol{B}) \times \boldsymbol{C}+\boldsymbol{B} \times(\boldsymbol{A} \times \boldsymbol{C})=\boldsymbol{A} \times(\boldsymbol{B} \times \boldsymbol{C})
$$

Using the above result, the curly bracket in the last term of Eq. (13.134) becomes $\boldsymbol{\alpha} \times\left(\boldsymbol{W}_{\mu} \times \boldsymbol{W}_{\nu}\right)$. Substituting it in Eq. (13.134), we get

$$
\begin{align*}
g\left(\boldsymbol{W}_{\mu}^{\prime} \times \boldsymbol{W}_{\nu}^{\prime}\right)= & g\left(\boldsymbol{W}_{\mu} \times \boldsymbol{W}_{\nu}\right)-g \boldsymbol{\alpha} \times\left(\boldsymbol{W}_{\mu} \times \boldsymbol{W}_{\nu}\right) \\
& -\boldsymbol{W}_{\mu} \times \partial_{\nu} \boldsymbol{\alpha}-\partial_{\mu} \boldsymbol{\alpha} \times \boldsymbol{W}_{\nu} \tag{13.135}
\end{align*}
$$

Using Eqs. (13.133) and (13.135), we get

$$
\begin{align*}
\mathcal{W}_{\mu \nu}^{\prime}-g\left(\boldsymbol{W}_{\mu}^{\prime} \times \boldsymbol{W}_{\nu}^{\prime}\right)= & \mathcal{W}_{\mu \nu}-g\left(\boldsymbol{W}_{\mu} \times \boldsymbol{W}_{\nu}\right) \\
& -\boldsymbol{\alpha} \times\left\{\boldsymbol{\mathcal { W }}_{\mu \nu}-g\left(\boldsymbol{W}_{\mu} \times \boldsymbol{W}_{\nu}\right)\right\} \tag{13.136}
\end{align*}
$$

Thus we find that

$$
\boldsymbol{W}_{\mu \nu}=\boldsymbol{\mathcal { W }}_{\mu \nu}-g\left(\boldsymbol{W}_{\mu} \times \boldsymbol{W}_{\nu}\right)=\partial_{\mu} \boldsymbol{W}_{\nu}-\partial_{\nu} \boldsymbol{W}_{\mu}-g\left(\boldsymbol{W}_{\mu} \times \boldsymbol{W}_{\nu}\right)
$$

is an iso-vector and the $\mathrm{SU}(2)$ gauge transformation is equivalent to rotating the vector in the iso-spin space (vide Problem 13.4.)
13.6 For invariance of Lagrangian under local gauge transformation under $\mathrm{SU}(2)$ group, three gauge vector fields $\boldsymbol{W}_{\mu}$ are introduced. The free Lagrangian density of the three gauge fields are

$$
\mathscr{L}_{g}=-\frac{1}{4} \boldsymbol{W}_{\mu \nu}^{a} \boldsymbol{W}_{a}^{\mu \nu}, \quad a=1,2,3
$$

We have shown in Solved Problem (13.5) that

$$
\boldsymbol{W}_{\mu \nu}=\partial_{\mu} \boldsymbol{W}_{\nu}-\partial_{\nu} \boldsymbol{W}_{\mu}-g\left(\boldsymbol{W}_{\mu} \times \boldsymbol{W}_{\nu}\right)
$$

is an iso-vector and the $\mathrm{SU}(2)$ local gauge transformation is equivalent to a rotation of the vector in the iso-spin space. Hence the scalar product

$$
\boldsymbol{W}_{\mu \nu} \cdot \boldsymbol{W}_{\mu \nu}=\sum_{a} \boldsymbol{W}_{\mu \nu}^{a} \boldsymbol{W}_{a}^{\mu \nu},
$$

is invariant under $\mathrm{SU}(2)$ local gauge transformation. Consequently, $\mathscr{L}_{g}$ is invariant under local $\mathrm{SU}(2)$ gauge transformation.
13.7 From the given relation $G_{F} / \sqrt{2}=1 /\left(2 v^{2}\right)$, we get

$$
\begin{aligned}
v^{2} & =\frac{1}{\sqrt{2} G_{F}}=6.0644 \times 10^{4} \\
v & =246.26 \mathrm{GeV}
\end{aligned}
$$

Given the fine structure constant, $e^{2} / 4 \pi=1 / 137$, we get the value of the coupling constant $e$.

$$
e=\left(\frac{4 \pi}{137}\right)^{1 / 2}=0.30286
$$

From Eq. (13.97), we get

$$
\begin{aligned}
g & =\frac{e}{\sin \theta_{W}}=\frac{0.30286}{\sqrt{0.225}}=0.6385 \\
g^{\prime} & =\frac{e}{\cos \theta_{W}}=\frac{e}{\sqrt{1-\sin ^{2} \theta_{W}}}=0.34403
\end{aligned}
$$

From Eqs. (13.93) and (13.100), we get the masses of the gauge bosons $M_{W}$ and $M_{Z}$.

$$
\begin{aligned}
M_{W} & =\frac{1}{2} v g=\frac{1}{2} \times 246.26 \times 0.6385=78.619 \mathrm{GeV} \\
M_{Z} & =\frac{M_{W}}{\cos \theta_{W}}=89.305 \mathrm{GeV}
\end{aligned}
$$

## Chapter 14

## Recent Developments

The purpose of the present book is limited to the study of Relativistic Quantum Mechanics and its extension to Feynman's formulation of Quantum Electrodynamics and highlight how Feynman's formulation is equivalent to the more general approach of Quantum Field Theory, the elements of which have been discussed in some detail. But, it is found desirable to give a bird's eye view of the difficulties encountered in treating higher order terms in the perturbation theory and the recent advances in Quantum Field Theory. The divergent integrals that one meets in the perturbation theory have been overcome by evolving renormalization techniques. It is shown that Quantum Electrodynamics is a gauge theory that is renormalizable. Recent advances in Quantum Field Theory include the successful attempts to bring in its fold the weak interactions and strong interactions. So, the Quantum Field Theory has emerged as a fundamental theory of electromagnetic, weak and strong interactions, leading to the formulation of the Standard Model of elementary particles. It is being generally felt that the Quantum Field Theory is a low energy limit of a more fundamental theory which includes gravitational interaction. The String Theory and Superstring Theory appear to be promising candidates for such a general theory of everything.

### 14.1 The Renormalization Program

As long as we are considering the QED processes in the lowest order of perturbation expansion, we face no problems and we get results that tally
fairly well with the experimental data. If we wish to improve the accuracy of our calculations and go to higher orders, we encounter divergent integrals as explained in Chap. 7. The integral ${ }^{1}$

$$
\int \frac{d^{4} k}{(2 \pi)^{4}} \frac{1}{\mathbf{k}^{4}}
$$

diverges logarithmically. Many have struggled with these infinities during 1940s. Eventually, they have found a way out by using different techniques of regularization and renormalization. By renormalizing mass, charge and wave functions, one can get rid of infinities in all orders of perturbation theory. A field theory which gives finite answers for all observable phenomena to all orders in perturbation theory by renormalization procedure is said to be renormalizable. QED is indeed a renormalizable theory.

### 14.1.1 Regularization and Renormalization

In evaluating the integral $\int \frac{d^{4} k}{(2 \pi)^{4}}$, we simply integrate only upto $\Lambda$, known as a cut off, to avoid the divergences. Then the integral is said to be regularized. After regularization, we speak of cut-off dependent quantities instead of divergent quantities.

But any physically observable quantity such as scattering cross section should not depend upon $\Lambda$.

In the S-matrix formalism, when the interaction is switched off, we deal with only bare particles without self-interaction i.e. with only bare mass $m_{0}$ and bare charge $e_{0}$. But these bare particles are not observable and all the observable attributes are experimentally observable mass $m$ and charge $e$. The difference between the bare mass $m_{0}$ and experimental mass $m$ and similarly the difference between bare charge $e_{0}$ and the observable charge $e$ depend upon the cut-off parameter $\Lambda$ but the physically observable quantities such as the scattering cross section depend only upon quantities such as experimental mass and charge. This is what we call renormalization. The physically observable quantities are indeed independent of the cut-off parameter. There are two different ways of achieving this. One is called the Pauli-Villars regularization and the other is known as the dimensional regularization.

[^80]
### 14.2 The Standard Model

It has been established that there are four different types of forces that operate in this universe: (1) Strong interaction (2) Electro-magnetic interaction (3) Weak interaction and (4) Gravitational interaction. It is the ultimate aim of the physicists to formulate a unified theory of all these four forces into one integrated "theory of everything".

Of the four interactions, electromagnetic interaction has been fairly well understood and all the processes in Quantum Electrodynamics can be explained to a high degree of accuracy on the basis of the theory developed by Tomonago, Schwinger and Feynman. The electro-magnetic interaction is mediated by photons (gauge Bosons) and in the field theoretic parlance, it is a gauge theory with gauge group $\mathrm{U}(1)$. The weak interaction is mediated by intermediate vector bosons $W^{+}, W^{-}, Z^{0}$ and Sheldon Glashow, Steven Weinberg and Abdus Salam developed a unified theory of electromagnetic and weak interactions which has come to be known as electro-weak theory which is a gauge theory of gauge group $\operatorname{SU}(2) \times$ $\mathrm{U}(1)$. The electro-weak theory has been enlarged to include strong interaction which arises from exchange of gluons between the coloured quarks that constitute the hadrons. The extended theory unifies the electroweak interactions with strong interactions and it is popularly known as the standard model of particle physics. It is a gauge theory of the electroweak and strong interactions with the gauge group $\mathrm{SU}(3) \times \mathrm{SU}(2) \times$ $\mathrm{U}(1)$.

### 14.2.1 Basic elementary particles - Fermions

The Standard Model is built on the basis of 12 elementary particles of spin $1 / 2$ known as Fermions which obey the Pauli exclusion principle. Each Fermion has a corresponding anti-particle. They are classified into quarks and leptons, depending upon their interactions. There are six quarks (up, down, strange, charm, bottom, top) and six leptons (electron, electron neutrino, muon, muon neutrino, tau, tau neutrino). Pairs of these particles from each group form a set, known as a generation, as shown in Table. 14.1 There are three generations.

The quarks carry colour charge and they interact by exchange of gluons which are also coloured objects. This results in the theory of colour forces known as quantum chromodynamics (QCD) and it has been developed as the theory of strong interactions in analogy with QED. But the fundamental difference is that photons do not carry electric charges in

Table 14.1: Fermions that form the basis for the Standard Model.

| Fermions | Charge | Generation I |  | Generation II | Generation III |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| Quarks | $+\frac{2}{3}$ | Up | $u$ | Charm | $c$ | Top | $t$ |
|  | $-\frac{1}{3}$ | Down | $d$ | Strange | $s$ | Bottom | $b$ |
| Leptons | -1 | Electron | $e^{-}$ | Muon | $\mu^{-}$ | Tau | $\tau^{-}$ |
|  | 0 | $e$-neutrino | $\nu_{e}$ | $\mu$-neutrino | $\nu_{\mu}$ | $\tau$-neutrino | $\nu_{\tau}$ |

QED but in QCD the gluons carry colour charges. QCD is a non-abelian gauge field theory of strong interactions and it has two important features:

1. At short distances i.e. at high momentum transfers, the interaction becomes weaker. This is known as the asymptotic freedom. In this region, one can resort to perturbation theory.
2. At large distances, the interaction grows stronger and stronger, leading to permanent confinement of quarks within colour-neutral composite particles (hadrons) consisting of either three quarks (baryons and hyperons) or either a quark-antiquark pair (mesons). This is known as the infra-red slavery.

Quarks also carry electric charge and weak isospin and hence they interact with other Fermions both by electromagnetic and weak interactions.

The remaining six Fermions do not carry colour charge and are called leptons. The three leptons, electron, muon and tau, are charged particles and they can interact through both electro-magnetic and weak interactions. The three neutrinos do not carry electric charge and hence they can interact only through weak interactions. This makes the detection of neutrinos an extremely difficult job.

Each member of a generation has greater mass than the corresponding particles of lower generations. The first generation charged particles do not decay. The second and third generation charged particles decay with very short half lives and are observed only in very high energy experi-
ments. Neutrinos of all generations do not decay but pervade the entire universe without much interaction with baryonic matter.

### 14.2.2 Force mediating particles - Bosons

The standard model explains the interaction between particles to arise from an exchange of Bosons which are the force-mediating particles.

1. Photons mediate the electromagnetic interactions between charged particles. The photon is massless and the electromagnetic interaction is well-described by the theory of quantum electrodynamics.
2. The intermediate vector bosons $W^{+}, W^{-}, Z^{0}$ mediate the weak interaction between leptons and quarks of different flavours. They are massive, $Z^{0}$ being more massive than $W^{ \pm}$. The weak interaction involving $W^{ \pm}$act exclusively on left-handed particles and righthanded antiparticles. Further, $W^{ \pm}$carry an electric charge of $\pm 1$ and they couple to the electromagnetic interaction. The electrically neutral $Z$ boson interacts with both left-handed particles and antiparticles. These three gauge bosons along with photon are clubbed together and they collectively mediate the electro-weak interactions.
3. The strong interaction between coloured quarks are mediated by eight gluons. The gluons are massless but they carry a colour which is labeled by a combination of colour and anti-colour charge (e.g. red-antigreen). Since the gluons have an effective colour charge, they can interact among themselves. The gluons and their interactions with themselves and with quarks are described by the theory of quantum chromodynamics.

### 14.2.3 The Higgs Boson

The Higgs particle is a massive scalar particle predicted by the Standard Model. It has no intrinsic spin and hence it is a Boson. The Higgs Boson plays a unique role in explaining how the elementary particles other than the photon and gluons acquire masses. It explains why the photon has no mass while the $W$ and $Z$ Bosons are very heavy. In electro-weak theory, the Higgs Boson generates the masses of leptons (electron, muon and tau particle) and quarks. It is the only fundamental particle predicted by the Standard Model and hence its discovery will validate the theoretical framework underlying the Standard Model. So, a massive effort was initiated by CERN in building up the Large Hadron Collder and sustained
the experimental activity, which finally succeeded in detecting this elusive particle in 2013.

### 14.2.4 Theoretical framework

The quantum field theory provides the theoretical framework for the Standard Model. Each particle is described in terms of a dynamical field that pervades the entire space-time. We first postulate a set of symmetries for the system and then write down the most general renormalizable Lagrangian for the system that observes these symmetries.

The global Poincaré symmetry is postulated for all relativistic quantum field theories. This symmetry includes the familiar translational symmetry, rotational symmetry and invariance under Lorentz transformation. The local $\mathrm{SU}(3) \times \mathrm{SU}(2) \times \mathrm{U}(1)$ gauge symmetry is an internal symmetry that essentially defines the standard model. The fields fall into different representations of the various symmetry groups of the Standard Model. The most general Lagrangian that displays the dynamics of the system involves 19 parameters whose numerical values are established by experiments.

Besides the symmetries postulated in the construction, the Standard Model exhibits four additional global symmetries, which are collectively denoted as accidental symmetries. These are continuous $\mathrm{U}(1)$ global symmetries which leave the Lagrangian invariant.

By Noether's theorem, each symmetry has an associated conservation law: the conservation of baryon number, electron number, muon number and tauon number. Each quark is assigned a baryon number $1 / 3$ while each antiquark is assigned a baryon number $-1 / 3$. Conservation of baryon number implies that the number of quarks minus the number of antiquarks is a constant. This is strongly supported by experiments and no violation has been observed for this conservation law.

Similarly, each electron and its associated neutrino is assigned an electron number +1 whereas each of its antiparticle is assigned an electron number -1 . In the same way, $\mu^{-}$and $\nu_{\mu}$ are assigned a muon number +1 whereas $\mu^{+}$and $\bar{\nu}_{\mu}$ are assigned a muon number $-1 ; \tau^{-}$and $\nu_{\tau}$ are assigned a tauon number +1 whereas $\tau^{+}$and $\bar{\nu}_{\tau}$ are assigned a tuaon number -1 . The standard model predicts that each of these lepton family numbers should be conserved separately in a way similar to the conservation of baryon number.

The standard model predicts that neutrinos are massless. But experiments indicate that neutrinos may have small masses, since neutrinos
are found to oscillate between flavours, signaling that the conservation of lepton family number is likely to be violated.

In addition to the accidental symmetries described above, the standard model exhibits several approximate symmetries.

### 14.2.5 Successes and failures of the Standard Model

The Standard Model predicted the existence of $W$ and $Z$ bosons, gluons, and the charm and top quarks before they were observed. Their predicted properties were experimentally confirmed with great precision as shown in Table. 14.2. The Standard Model also made several predictions about the decay of $Z$ bosons, which were confirmed by the Large Electron-Positron Collider experiments at CERN.

Table 14.2: Predictions of the masses of $W$ and $Z$ bosons by the Standard Model (SM) and their experimentally measured values.

| Quantity | SM prediction $(\mathrm{GeV})$ | Experiment $(\mathrm{GeV})$ |
| :---: | :---: | :---: |
| Mass of $W$ Boson | $80.390 \pm 0.018$ | $80.398 \pm 0.025$ |
| Mass of $Z$ Boson | $91.1874 \pm 0.0021$ | $91.1876 \pm 0.0021$ |

There is some experimental evidence for the oscillations of neutrino from one flavour to the other. This will be possible only if the neutrinos possess some mass as against the zero mass prediction by the standard model. If the neutrinos have small masses, then there will be some violation of the lepton number flavour conservation.

One of the predictions of the Standard Model is the Higgs Boson which has been experimentally observed in CERN. The Standard Model involves 19 numerical constants whose values are unrelated and arbitrary. and it does not allow finite, although small, masses for neutrinos.

Gravitation is out of the purview of the Standard Model and there is no way of including the general relativity. So, the Standard Model does not explain the initial conditions of the universe that give rise to certain observed properties of the present-day universe, properties such as the predominance of matter over antimatter (matter/antimatter asymmetry), and its isotropy and homogeneity over large distances.

### 14.3 Gravitational Interaction

We have seen that the Standard Model explains fairly well the world of elementary particles, on the basis of electromagnetic, weak and strong interactions. The only other force left out is gravitational interaction. That means the gravitational interaction is so weak compared to electromagnetic, weak and strong interactions at the present day energies that the accelerators can reach. To get a feeling of the relative strength, consider the ratio of gravitational force to the electromagnetic force that operates between two protons:

$$
\begin{equation*}
\frac{F_{\text {gravity }}}{F_{\text {e.m. }}}=\frac{G M_{p}^{2} / r^{2}}{e^{2} / r^{2}} \approx 10^{-36} \ll 1 \tag{14.1}
\end{equation*}
$$

where $G$ is the gravitational constant and $M_{p}$ is the rest mass of the proton. At high energies, one has to replace the rest mass of the proton by the effective relativistic mass

$$
\begin{equation*}
M_{\mathrm{eff}}=\frac{M_{p}}{\left(1-\frac{v^{2}}{c^{2}}\right)^{1 / 2}} . \tag{14.2}
\end{equation*}
$$

At high energies, when $M_{\text {eff }}=10^{18} M_{p}$, the gravitational force is comparable to electromagnetic force and hence cannot be ignored. At such high energies, the Standard Model should fail.

From the general theory of relativity, it is conjectured that the gravitational field arises from an exchange of gravitons which are massless but have spin 2. This has motivated physicists to look for a more fundamental theory which will unify all the four forces that operate between the elementary particles at high energies that will reduce to the Standard Model at low energies.

Einstein's general theory of relativity relates the gravitational forces to the curved nature of space-time. This interlinking of gravity with space-time offers insurmountable difficulties in finding a quantum theory of gravity that will unify all the four fundamental forces. Such a quantum theory should describe our universe, the physics of black holes and the processes involving elementary particles. String theory appears to be a promising candidate for such a unified theory.

### 14.4 Birth of String Theory

In the late 1960s, a totally novel idea emerged that particles should be treated as string-like objects instead of point-like objects under the name
"dual models". This was to explain the experimentally observed behaviour of strongly interacting particles, called hadrons, at large angular momentum $J$. The experimental data on resonances (excited states of the light hadrons) of mass $M$ and spin $J$ fitted well with the relation $M^{2}=J \alpha^{\prime}$, where $\alpha^{\prime}$ is called the Regge slope. These strings have infinitesimal thickness and tension $T=1 /\left(2 \pi \alpha^{\prime}\right)$ and can be either open or closed. The particles can be identified as the various vibrational modes of an open string or closed string ${ }^{2}$. Since the vibrational modes predicted massless particles of spin 2 and higher spins which are not seen in the hadronic world, the theory was discorded.

In 1974, Scherk and Schwarz ${ }^{3}$ made an important observation that the spin-2 massless particle in the spectrum of vibrational modes of the string should be identified as the graviton. It has been known to the theoretical physicists for a long time that if a good quantum theory of gravity is to be developed then the particle that would carry the gravitational force would have zero mass and two units of spin. Thus the theory of "dual models" was revived under the name "string theory" in the hope of formulating a quantum theory of gravity consistent with the laws of quantum mechanics.

### 14.4.1 Energy and Length Scales

Hitherto, there is no direct experimental observation that supports the string-like nature of the elementary particles. It is because the presentday accelerators have not reached the energy scale at which the string-like nature can be detected. For the applicability of String Theory, we need to discuss the Planckian units to determine the energy scale.

In Relativistic Quantum Mechanics and Quantum Field Theory, we use natural units with $c=1$ and $\hbar=1$. In String Theory, one has to use Planckian units with Newton's Gravitational constant $G=1$ besides $\hbar=1$ and $c=1$. This yields Planck length $l_{\mathrm{pl}}$, Planck mass $m_{\mathrm{pl}}$ and

[^81]Planck time $t_{\mathrm{pl}}$ in terms of $G, \hbar$ and $c$.

$$
\begin{align*}
l_{\mathrm{pl}} & =\sqrt{\frac{G \hbar}{c^{3}}}=1.61 \times 10^{-33} \mathrm{~cm} \\
m_{\mathrm{pl}} & =\sqrt{\frac{\hbar c}{G}}=2.17 \times 10^{-5} \mathrm{gm}  \tag{14.3}\\
t_{\mathrm{pl}} & =\frac{l_{\mathrm{pl}}}{c}=5.4 \times 10^{-44} \mathrm{sec}
\end{align*}
$$

Using Einstein's mass-energy relation $E=m c^{2}$, the Planck mass $m_{\mathrm{pl}}=$ $10^{19} \mathrm{GeV}$, in energy units. Our present day accelerators have reached the energy only upto $10 \mathrm{TeV}=10^{4} \mathrm{GeV}=10^{13} \mathrm{eV}$ and the Planck energy is far above. This means that to detect the string-like nature, we need the length scale of $10^{-33} \mathrm{~cm}$ or the energy scale of $10^{19} \mathrm{GeV}$ for the accelerators.

It is expected that if the String Theory is a more fundamental unified theory, constructed at $10^{19} \mathrm{GeV}$, then it must reduce to the Standard Model at 100 GeV .

### 14.4.2 Basic theory of strings

The basic concepts used in particle physics are extended to the theory of strings ${ }^{4}$. The trajectory of a particle in classical mechanics becomes a world line in relativistic mechanics. Characterizing the world line by a parameter $\tau$, the coordinates of the points on the world line can be represented by $x^{\mu}(\tau)=\boldsymbol{x}(\tau), t(\tau)$. In Fig. 14.1, we plot a world line on a $2+1$ dimensional space-time.

For more than two dimensional space, it is not possible to draw the world line on the plane of a paper but it can be visualized. Given the two space-time points A and B , one can draw very many trajectories connecting the points A and B , but the chosen path is the one which extremises the interval between A and B with the constraint that $\delta\left(x^{\mu}\right)=0$ at the end points A and B. This will yield a set of second order differential equations for $x^{\mu}(\tau)$, the solution of which describes the classical path.

### 14.4.3 Open strings

If we replace the point-like particles by open strings of length $l$, the trajectories of the string will sweep out a two-dimensional surface called world

[^82]

Figure 14.1: World line of a particle parametrized by $\tau$
sheet that describes the time evolution of the string. The world sheet can be parametrized by two parameters $\sigma, \tau$, where $0 \leq \sigma \geq l$ represents the points on the string of length $l$. The world sheet can be described by the coordinates $x^{\mu}(\sigma, \tau)$, where the superscript $\mu=0,1,2, \cdots d$ denotes the components of $d+1$ dimensional space-time coordinate system. In particle physics, we have always assumed $3+1$ dimensional space-time but in string theory, the space-time dimensions are determined by the stringent conditions imposed by quantum mechanics and relativity. Such a $d+1$ dimensional space-time is called the target space in which the string generates a world sheet.

The open string has two end points $\sigma=0$ and $\sigma=l$, called the boundary points. To find the classical trajectory of particles, we have extremised the interval between the space-time points A and B. In the case of strings, we need to extremise the area of the world sheet surface to get a set of second order differential equations for $x^{\mu}(\sigma, \tau)$. In particle theory, we have used the boundary conditions $\delta\left(x^{\mu}\right)=0$ at A and B. In the case of open strings, we use either the Neumann boundary conditions at the boundary points

$$
\begin{equation*}
\left.\frac{\partial}{\partial \sigma} x^{\mu}(\sigma, \tau)\right|_{\sigma=0}=\left.\frac{\partial}{\partial \sigma} x^{\mu}(\sigma, \tau)\right|_{\sigma=l}=0 \tag{14.4}
\end{equation*}
$$

at some directions $\mu=0,1,2, \cdots, p$ where $p \leq d$ and Dirichlet boundary conditions

$$
\begin{equation*}
x^{m}(\sigma=0, \tau)=b^{m}, \quad x^{m}(\sigma=l, \tau)=c^{m} \tag{14.5}
\end{equation*}
$$

(where $b_{m}$ and $c_{m}$ are constants) for the remaining directions $m=p+$ $1, \cdots, d)$. There can also be mixed boundary conditions where some directions have Neumann boundary conditions at $\sigma=0$ and Dirichlet boundary conditions at $\sigma=l$.

Initially, the Dirichlet boundary conditions were considered as unphysical since they were not consistent with the laws of the special theory of relativity in $d+1$ dimensions and hence strings with only the Newmann boundary conditions were studied. Later, it was pointed out that Nature requires that the laws of special theory of relativity need be satisfied only in $3+1$ dimensional space and hence Dirichlet boundary conditions are allowed in other directions. It is found that the Dirichlet boundary conditions lead to non-perturbative spectrum of string theory.

Only the perturbative spectrum involves fundamental strings where the boundary conditions are only Neumann conditions. Extremising the surface area of the world sheet of the open string subjected to the boundary conditions, we obtain the classical solution $x^{\mu}(\sigma, \tau)$, which can be expanded into Fourier modes $a_{n}^{\mu}$ that resembles the linear superposition of standing waves. Considering $x^{\mu}(\sigma, \tau)$ as an operator, quantization results in commutation relations between the Fourier modes $a_{n}^{\mu}$. The modes with positive $n$ resemble harmonic oscillator annihilation operators and those with negative $n$ resemble creation operators $a_{-n}^{\mu}=\left(a_{n}^{\mu}\right)^{\dagger}$.

The non-negative norm imposed by quantum mechanics on the string states and the requirement to satisfy the symmetries of the special theory of relativity results in the $25+1$ dimensions for space-time. Thus the space-time dimension emerges from the string theory instead of being assumed as $3+1$ in quantum field theory.

The first excited string state in $25+1$ dimensional space-time is identified as the massless photon. The vacuum state has negative mass-squared which is usually referred to as tachyons.

### 14.4.4 Closed strings

Unlike open strings, closed strings have no boundary points. The point $\sigma$ on the closed string coincides with the point $\sigma+l$ since the ends of the string of length $l$ are joined to form a closed loop.

$$
\begin{equation*}
x^{\mu}(\sigma, \tau)=x^{\mu}(\sigma+l, \tau) \tag{14.6}
\end{equation*}
$$

Again, the classical solution is found by extremising the area of the closed string world sheet subject to the above condition.

It is found that the first excited massless state can be identified as the graviton which is the quantum responsible for gravitational force. Thus the closed string theory incorporates the quantum of gravity. The ground state in closed string theory is again a negative mass-squared state which represents closed string tachyons.

For the study of interactions and scattering, perturbative string theory is developed in a way analogous to perturbation expansion in particle theory using Feynman graphs. In perturbative string theory, Feynmanlike diagrams that are used are called fat graphs.

### 14.5 Superstring Theories

The string theory that is developed in $25+1$ dimensional space is called Bosonic String Theory since all the vibrational modes of the string correspond to integer spins. Generally Bosons represent particles that transmit forces and the Fermions are particles that make up the matter. If you impose supersymmetry, then for every Bosonic state, there is a corresponding Fermion and the string theory so obtained is called Superstring Theory. There are five kinds of superstring theories, as shown in the table below.

The most important point in formulating a string theory is that it should obey quantum mechanics in a sensible way. For bosonic strings, this is satisfied only if the space-time dimensions number 26 . For superstrings we can whittle it down to 10 . How we get down to the four space-time dimensions, that we live in, is a different story.

If we ask how to get from ten space-time dimensions to four spacetime dimensions, there are many possible ways by which the superfluous six dimensions are made very much smaller than the other four in string theory. This process of compactification of unwanted space-time dimensions yields interesting physics on its own.

But the number of string theories has also been shrinking in recent years, because string theorists are discovering that what they thought were completely different theories were in fact different ways of looking at the same theory!

This period in string history has been given the name the second string revolution. And now the biggest question is how to collapse the table above into one theory, which some people want to call the $M$ theory, for it is the Mother of all theories.

At one time, string theorists believed there were five distinct superstring theories: type I, types IIA and IIB, and the two heterotic string

Table 14.3: Table of String Theories

| Type | Space-time <br> Dimensions | Details |
| :---: | :---: | :---: |
| Bosonic | 26 | Only bosons, no fermions means only forces, no matter, with both open and closed strings. Major flaw: a particle with imaginary mass, called the tachyon |
| I | 10 | Supersymmetry between forces and matter, with both open and closed strings, no tachyon, group symmetry is $\mathrm{SO}(32)$ |
| IIA | 10 | Supersymmetry between forces and matter, with closed strings only, no tachyon, massless fermions spin both ways (nonchiral) |
| IIB | 10 | Supersymmetry between forces and matter, with closed strings only, no tachyon, massless fermions only spin one way (chiral) |
| HO | 10 | Supersymmetry between forces and matter, with closed strings only, no tachyon, heterotic, meaning right moving and left moving strings differ, group symmetry is $\mathrm{SO}(32)$ |
| HE | 10 | Supersymmetry between forces and matter, with closed strings only, no tachyon, heterotic, meaning right moving and left moving strings differ, group symmetry is E8 $\times$ E8 |

theories.
But now it is known that this naive picture was wrong, and that the five superstring theories are connected to one another as if they are each a special case of a more fundamental M theory. These theories are related by transformations that are called dualities. If two theories are related by a duality transformation, it means that the first theory can be transformed in some way so that it ends up looking just like the second theory. The two theories are then said to be dual to one another under that kind of transformation.

These dualities link quantities that were also thought to be separate. Large and small distance scales, strong and weak coupling strengths -
these quantities have always marked very distinct limits of behavior of a physical system, in both classical field theory and quantum particle physics. But strings can obscure the difference between large and small, strong and weak, and this is how these five very different theories end up being related.

## Large and small distance

The duality symmetry that obscures our ability to distinguish between large and small distance scales is called T-duality, and comes about from the compactification of extra space dimensions in a ten dimensional superstring theory.

Suppose we are in ten space-time dimensions, which means we have nine space and one time. Take one of those nine space dimensions and make it a circle of radius $R$, so that traveling in that direction for a distance $L=2 \pi R$ takes you around the circle and brings you back to where you started.

A particle traveling around this circle will have a quantized momentum around the circle, and this will contribute to the total energy of the particle. But a string is very different, because in addition to traveling around the circle, the string can wrap around the circle. The number of times the string winds around the circle is called the winding number, and that is also quantized.

Now the weird thing about string theory is that these momentum modes and the winding modes can be interchanged, as long as we also interchange the radius $R$ of the circle with the quantity $2 L_{\mathrm{st}} / R$, where $L_{\mathrm{st}}$ is the string length. If $R$ is very much smaller than the string length, then the quantity $2 L_{\mathrm{st}} / R$ is going to be very large. So exchanging momentum and winding modes of the string exchanges a large distance scale with a small distance scale.

This type of duality is called T-duality. T-duality relates Type IIA superstring theory to Type IIB superstring theory. That means if we take Type IIA and Type IIB theory and compactify them both on a circle, then switching the momentum and winding modes, and switching the distance scale, changes one theory into the other! The same is also true for the two heterotic theories.

So T-duality obscures the difference between large and small distances. What looks like a very large distance to a momentum mode of a string looks like a very small distance to a winding mode of a string.

## Strong and weak coupling

What is a coupling constant? This is some number that tells us how strong an interaction is. A larger coupling constant means a stronger force, and a smaller coupling constant means a weaker force.

For QED, the coupling constant is proportional to the square of the electric charge and the perturbation theory is usually employed to solve any problem. But if the coupling constant is larger, this method of calculation breaks down.

This also can happen in string theory. String theories have a coupling constant. But unlike in particle theories, the string coupling constant is not just a number, but depends on one of the oscillation modes of the string, called the dilaton. Exchanging the dilaton field with minus itself exchanges a very large coupling constant with a very small one. This symmetry is called S-duality. If two string theories are related by S-duality, then one theory with a strong coupling constant is the same as the other theory with weak coupling constant. Since the theory with strong coupling cannot be understood by means of expanding in a series, but the theory with weak coupling can. So if the two theories are related by S-duality, then we just need to understand the weak theory, and that is equivalent to understanding the strong theory.

Superstring theories related by S-duality are: Type I superstring theory with heterotic $\mathrm{SO}(32)$ superstring theory, and Type IIB theory with itself.

Starting from any dimension greater than four, it is necessary to consider how these are reduced to four dimensional space-time.

Two different ways have been proposed to resolve this apparent contradiction. The first is to compactify the extra dimensions; i.e., the 6 or 7 extra dimensions are so small as to be undetectable by present day experiments. A standard analogy for this is to consider multidimensional space as a garden hose. If the hose is viewed from a sufficient distance, it appears to have only one dimension, its length. Indeed, think of a ball just small enough to enter the hose. Throwing such a ball inside the hose, the ball would move more or less in one dimension; in any experiment we make by throwing such balls in the hose, the only important movement will be one-dimensional, that is, along the hose. However, as one approaches the hose, one discovers that it contains a second dimension, its circumference. Thus, an ant crawling inside it would move in two dimensions (and a fly flying in it would move in three dimensions). This "extra dimension" is only visible within a relatively close range to
the hose, or if one "throws in" small enough objects. Similarly, the extra compact dimensions are only "visible" at extremely small distances, or by experimenting with particles with extremely small wavelengths (of the order of the compact dimension's radius), which in quantum mechanics means very high energies (see wave-particle duality).

T-duality relates the large and small distance scales between string theories, whereas S-duality relates strong and weak coupling strengths between string theories. U-duality links T-duality and S-duality. One such theory is the 11 -dimensional M-theory, which requires space-time to have eleven dimensions, as opposed to the usual three spatial dimensions and the fourth dimension of time. The original string theories from the 1980s describe special cases of M-theory where the eleventh dimension is a very small circle or a line, and if these formulations are considered as fundamental, then string theory requires ten dimensions. But the theory also describes universes like ours, with four observable space-time dimensions, as well as universes with up to 10 flat space dimensions, and also cases where the position in some of the dimensions is not described by a real number, but by a completely different type of mathematical quantity. So the notion of space-time dimension is not fixed in string theory: it is best thought of as different in different circumstances.

## Appendix A

## List of Symbols and Notation

$\sigma_{x}, \sigma_{y}, \sigma_{z}$
Pauli spin matrices
$\begin{aligned} \sigma_{x} & =\left[\begin{array}{rr}0 & 1 \\ 1 & 0\end{array}\right], \quad \sigma_{y}=\left[\begin{array}{rr}0 & -i \\ i & 0\end{array}\right], \\ \sigma_{z} & =\left[\begin{array}{rr}1 & 0 \\ 0 & -1\end{array}\right] .\end{aligned}$
$\beta, \alpha_{x}, \alpha_{y}, \alpha_{z}$
$\gamma_{0}, \gamma_{x}, \gamma_{y}, \gamma_{z}$
$\gamma_{5}, \gamma_{5}^{\prime}$
Dirac matrices
$\beta=\left[\begin{array}{rr}I & 0 \\ 0 & -I\end{array}\right], \quad \boldsymbol{\alpha} \cdot \boldsymbol{p}=\left[\begin{array}{ll}0 & \boldsymbol{\sigma} \\ \boldsymbol{\sigma} & 0\end{array}\right]$.
Gamma matrices
$\gamma_{0}=\beta, \quad \gamma_{x}=\beta \alpha_{x}, \quad \gamma_{y}=\beta \alpha_{y}, \quad \gamma_{z}=\beta \alpha_{z}$.
$\gamma_{0}=\left[\begin{array}{rr}I & 0 \\ 0 & -I\end{array}\right], \quad \gamma=\left[\begin{array}{rr}0 & \boldsymbol{\sigma} \\ -\boldsymbol{\sigma} & 0\end{array}\right]$.
$\gamma_{5}=\gamma_{0} \gamma_{x} \gamma_{y} \gamma_{z}, \quad \gamma_{5}^{\prime}=i \gamma_{5}$
$\gamma_{5}=-i\left[\begin{array}{cc}0 & I \\ I & 0\end{array}\right], \quad \gamma_{5}^{\prime}=\left[\begin{array}{cc}0 & I \\ I & 0\end{array}\right]$.
Scalar product of three-vectors

Scalar product of four-vectors

Contravariant
four-vector
Covariant
four-vector
$\boldsymbol{p} \cdot \boldsymbol{x}=p_{x} x+p_{y} y+p_{z} z ;$
$\boldsymbol{A} \cdot \boldsymbol{B}=A_{x} B_{x}+A_{y} B_{y}+A_{z} B_{z}$.
$\mathbf{p} \cdot \mathbf{x}=E t-\boldsymbol{p} \cdot \boldsymbol{x} ; \quad \mathbf{A} \cdot \mathbf{B}=A_{0} B_{0}-\boldsymbol{A} \cdot \boldsymbol{B}$.
The upright bold letters denote four-vectors and the bold italics denote three-vectors.
$x^{\mu}, \mu=0,1,2,3$.
$x^{0}=c t, x^{1}=x, x^{2}=y, x^{3}=z$.
$x_{\mu}, \mu=0,1,2,3$.
$x_{0}=c t, x_{1}=-x, x_{2}=-y, x_{3}=-z$.

| $\not p$ | $\gamma^{\mu} p_{\mu}=\gamma^{0} p_{0}-\gamma \cdot \boldsymbol{p}=\gamma^{0} p_{0}-\gamma_{x} p_{x}-\gamma_{y} p_{y}-\gamma_{z} p_{z}$. |
| :---: | :---: |
| $g_{\mu \nu}$ | Components of metric tensor $=g^{\mu \nu}, \quad \operatorname{diag}(1,-1,-1,-1)$ |
| $\square$ | D'Alembertian operator $\square=\frac{\partial}{\partial x_{\mu}} \frac{\partial}{\partial x^{\mu}}=\frac{1}{c^{2}} \frac{\partial^{2}}{\partial t^{2}}-\nabla^{2} .$ |
| $\partial_{\mu}$ | $\frac{\partial}{\partial x^{\mu}} .$ |
| $\partial^{\mu}$ | $\frac{\partial}{\partial x_{\mu}} .$ |
| $F^{\mu \nu}$ | Electromagnetic field tensor $F^{\mu \nu}=\partial^{\nu} A^{\mu}-\partial^{\mu} A^{\nu}=\frac{\partial A^{\mu}}{\partial x_{\nu}}-\frac{\partial A^{\nu}}{\partial x_{\mu}}=-F^{\nu \mu} .$ |
| $[A, B]_{-}$ | $[A, B]_{-}=A B-B A . \quad$ (Commutator bracket) |
| $\{A, B\}_{+}$ | $\{A, B\}_{+}=A B+B A$. (Anticommutator bracket) |
| $\mathscr{L}$ | Lagrangian density. |
| $L$ | Total Lagrangian $L=\int d^{3} x \mathscr{L}$. |
| $\mathscr{A}$ | Action $\mathscr{A}=\int d t L$. |
| $\mathscr{H}$ | Hamiltonian density. |
| H | Total Hamiltonian $H=\int d^{3} x \mathscr{H}$. |

Consider two events at space-time points $(x, y, z, t)$ and $(x+d x, y+$ $d y, z+d z, t+d t)$. Then the interval $d s$ between these two space-time points is defined by

$$
\begin{equation*}
d s^{2}=c^{2} d t^{2}-\left(d x^{2}+d y^{2}+d z^{2}\right) \tag{A.1}
\end{equation*}
$$

The square of the interval $d s^{2}$ is the same for all observers in different inertial frames of reference and so is invariant under Lorentz transformation. With this definition, we say that the space-time interval is time-like if

$$
d s^{2}>0
$$

and space-like if

$$
d s^{2}<0
$$

and light-like if

$$
d s^{2}=0
$$

The invariant space-time interval $d s$ is obtained by defining two types of four-vectors, one contravariant four-vector $x^{\mu}$ with an upper index $\mu$ and the other called covariant four-vector $x_{\mu}$ with a lower index $\mu$ and taking a scalar product (inner product) of them.

## Contravariant four-vector

$$
\begin{equation*}
x^{\mu}=x^{0}, x^{1}, x^{2}, x^{3}=c t, x, y, z . \tag{A.2}
\end{equation*}
$$

## Covariant four-vector

$$
\begin{equation*}
x_{\mu}=x_{0}, x_{1}, x_{2}, x_{3}=c t,-x,-y,-z . \tag{A.3}
\end{equation*}
$$

## Scalar product

$$
\begin{equation*}
d s^{2}=\sum_{\mu=0}^{3} d x^{\mu} d x_{\mu}=c^{2} d t^{2}-d x^{2}-d y^{2}-d z^{2} \tag{A.4}
\end{equation*}
$$

The relation between the contravariant and covariant vectors can be expressed by defining a metric tensor $g_{\mu \nu}$

$$
g_{\mu \nu}=\left[\begin{array}{rrrr}
1 & 0 & 0 & 0  \tag{A.5}\\
0 & -1 & 0 & 0 \\
0 & 0 & -1 & 0 \\
0 & 0 & 0 & -1
\end{array}\right]
$$

Then

$$
\begin{align*}
x_{\mu} & =\sum_{\nu} g_{\mu \nu} x^{\nu}=g_{\mu \nu} x^{\nu} \\
& =g_{\mu 0} x^{0}+g_{\mu 1} x^{1}+g_{\mu 2} x^{2}+g_{\mu 3} x^{3} . \tag{A.6}
\end{align*}
$$

The repeated index $\nu$ implies summation over $\nu$ and if this convention is used, then the summation sign is redundant. The relation (A.6) yields

$$
\begin{equation*}
x_{0}=x^{0}, \quad x_{1}=-x^{1}, \quad x_{2}=-x^{2}, \quad x_{3}=-x^{3} . \tag{A.7}
\end{equation*}
$$

Since the geometry of the space-time depends on the distribution of matter, leading to the concept of curved space-time, the metric tensor $g_{i j}$ is a function of the space-time coordinates.

The contravariant metric tensor $g^{\mu \nu}$ is given by

$$
g^{\mu \nu}=g_{\mu \nu} .
$$

A contravariant four-vector $A^{\mu}$ transforms from one coordinate system (unprimed) $x^{\mu}$ to another coordinate system (primed) $x^{\mu}$ according to the formula

$$
\begin{equation*}
A^{\prime \mu}=\frac{\partial x^{\prime \mu}}{\partial x^{\nu}} A^{\nu} ; \quad A^{\nu}=\frac{\partial x^{\nu}}{\partial x^{\prime \mu}} A^{\prime \mu} \tag{A.8}
\end{equation*}
$$

A covariant four-vector $A_{\mu}$ transforms from the unprimed coordinate system $x_{\mu}$ to the primed coordinate system $x_{\mu}^{\prime}$ according to the law

$$
\begin{equation*}
A_{\mu}^{\prime}=\frac{\partial x^{\nu}}{\partial x^{\prime \mu}} A_{\nu} ; \quad A_{\nu}=\frac{\partial x^{\prime \mu}}{\partial x^{\nu}} A_{\mu}^{\prime} \tag{A.9}
\end{equation*}
$$

If $\left(A^{0}, A^{1}, A^{2}, A^{3}\right)$ are the contravariant components of a four-vector $\mathbf{A}$, then its covariant components $A_{\mu}$ are given by

$$
\begin{align*}
& A_{\mu}=g_{\mu \nu} A^{\nu}, \quad(\mu=0,1,2,3)  \tag{A.10}\\
& A_{0}=A^{0} ; \quad A_{1}=-A^{1} ; \quad A_{2}=-A^{2} ; \quad A_{3}=-A^{3} \tag{A.11}
\end{align*}
$$

Thus, we can write the contravariant and covariant components of a fourvector as

$$
A^{\mu}=\left(A^{0}, \boldsymbol{A}\right) ; \quad A_{\mu}=\left(A^{0},-\boldsymbol{A}\right) .
$$

The scalar product (inner product) of the two four-vectors $\mathbf{A}$ and $\mathbf{B}$ is

$$
\begin{align*}
\mathbf{A} \cdot \mathbf{B}=A^{\mu} B_{\mu} & =A^{0} B_{0}+A^{1} B_{1}+A^{2} B_{2}+A^{3} B_{3} \\
& =A_{0} B_{0}-A_{x} B_{x}-A_{y} B_{y}-A_{z} B_{z} \\
& =A^{0} B_{0}-\boldsymbol{A} \cdot \boldsymbol{B} . \tag{A.12}
\end{align*}
$$

We adopt the convention of representing four-vectors by upright bold letters and three-component vectors by italic bold letters.

The differential operators in the four-dimensional space-time $x^{\mu}$ transforms as the components of a covariant four-vector.

$$
\begin{equation*}
\frac{\partial}{\partial x^{\mu}}=\frac{\partial x^{\nu}}{\partial x^{\prime \mu}} \frac{\partial}{\partial x^{\nu}} \tag{A.13}
\end{equation*}
$$

Thus the differential operator is a covariant four-vector.

$$
\begin{equation*}
\partial_{\mu}=\frac{\partial}{\partial x^{\mu}}=\left(\frac{\partial}{\partial x^{0}}, \nabla\right) . \tag{A.14}
\end{equation*}
$$

The corresponding contravariant four-vector is

$$
\begin{equation*}
\partial^{\nu}=g^{\mu \nu} \partial_{\nu}=\frac{\partial}{\partial x_{\nu}}=\left(\frac{\partial}{\partial x_{0}},-\boldsymbol{\nabla}\right) . \tag{A.15}
\end{equation*}
$$

From the foregoing discussion, it is clear that $g_{\mu \nu}$ contains all the information about the geometry of the space - in this case, Minkowski space-time. If we confine ourselves to the special theory of relativity, the metric tensor $g_{\mu \nu}$ plays a passive role but it plays an active role in general relativity since the space-time geometry is not fixed in advance and can be curved depending on the distribution of matter. So, in special relativity, one can avoid the subtle distinction between the contravariant and covariant vectors and define the scalar product of four-vectors

$$
\begin{array}{r}
\mathbf{A} \cdot \mathbf{B}=A_{0} B_{0}-\boldsymbol{A} \cdot \boldsymbol{B}=A_{0} B_{0}-A_{x} B_{x}-A_{y} B_{y}-A_{z} B_{z} . \\
\quad \not p=\gamma_{\mu} p_{\mu}=\gamma_{0} p_{0}-\gamma \cdot \boldsymbol{p}=\gamma_{0} p_{0}-\gamma_{x} p_{x}-\gamma_{y} p_{y}-\gamma_{z} p_{z} . \tag{A.17}
\end{array}
$$

By the above definition, we can avoid altogether the concept of contravariant and covariant vectors as was done by Feynman in his book on Quantum Electrodynamics (1962). We have followed this convention in chapters 1-7 and only in chapters 8 -14, we have introduced the notation of contravariant and covariant vectors.

## Appendix B

## The Density of Final States

## 1. One-particle final state

Volume of phase space $=V d^{3} p=V p^{2} d p d \Omega$.
The number of states $d n$ available for a particle with momentum lying between $p$ and $p+d p$ is obtained by dividing the volume of phase space by the volume of the unit cell $h^{3}$.

$$
d n=\frac{V p^{2} d p d \Omega}{h^{3}}=\frac{V p^{2} d p d \Omega}{(2 \pi \hbar)^{3}}
$$

Let us assume the quantization volume to be unity i.e. $V=1$ and use natural units i.e. $\hbar=1, c=1$. The density of states $\rho$ is defined as the number of states available per unit energy interval for a particle scattered within a solid angle $d \Omega$.

$$
\begin{equation*}
\rho(E, p)=\frac{d n}{d E d \Omega}=\frac{p^{2}}{(2 \pi)^{3}} \frac{d p}{d E}=\frac{E p}{(2 \pi)^{3}} . \tag{B.1}
\end{equation*}
$$

The last step is obtained by using the relativistic relation between energy and momentum $E^{2}=p^{2}+m^{2}$ such that $d E / d p=p / E$.

## 2. Two-particle final state

Here we have two particles in the final state with total energy $E$ and total momentum $\boldsymbol{P}$ and the observation is made on the particle 1.

$$
\begin{align*}
& E=E_{1}+E_{2} ; \quad \boldsymbol{P}=\boldsymbol{p}_{1}+\boldsymbol{p}_{2} .  \tag{B.2}\\
& \rho\left(E ; p_{1}, p_{2}\right)=\frac{p_{1}^{2}}{(2 \pi)^{3}} \frac{d p_{1}}{d E} . \tag{B.3}
\end{align*}
$$

We need to evaluate $d E / d p_{1}$.

$$
\begin{equation*}
\frac{d E}{d p_{1}}=\frac{d E_{1}}{d p_{1}}+\frac{d E_{2}}{d p_{1}}=\frac{p_{1}}{E_{1}}+\frac{d E_{2}}{d p_{2}} \frac{d p_{2}}{d p_{1}}=\frac{p_{1}}{E_{1}}+\frac{p_{2}}{E_{2}} \frac{d p_{2}}{d p_{1}} . \tag{B.4}
\end{equation*}
$$

To evaluate $d p_{2} / d p_{1}$, we proceed as follows:

$$
\begin{align*}
& \boldsymbol{p}_{2}=\boldsymbol{P}-\boldsymbol{p}_{1} ; \quad p_{2}^{2}=P^{2}+p_{1}^{2}-2 \boldsymbol{P} \cdot \boldsymbol{p}_{1} ;  \tag{B.5}\\
& 2 p_{2} \frac{d p_{2}}{d p_{1}}=2 p_{1}-2 \frac{\boldsymbol{P} \cdot \boldsymbol{p}_{1}}{p_{1}} \\
& \frac{d p_{2}}{d p_{1}}=\frac{p_{1}}{p_{2}}-\frac{\boldsymbol{P} \cdot \boldsymbol{p}_{1}}{p_{1} p_{2}}=\frac{p_{1}^{2}-\boldsymbol{P} \cdot \boldsymbol{p}_{1}}{p_{1} p_{2}} . \tag{B.6}
\end{align*}
$$

Substituting (B.6) into (B.4), we get

$$
\begin{align*}
\frac{d E}{d p_{1}} & =\frac{p_{1}}{E_{1}}+\frac{p_{2}}{E_{2}} \frac{p_{1}^{2}-\boldsymbol{P} \cdot \boldsymbol{p}_{1}}{p_{1} p_{2}} \\
& =\frac{p_{1}^{2} E_{2}+\left(p_{1}^{2}-\boldsymbol{P} \cdot \boldsymbol{p}_{1}\right) E_{1}}{E_{1} E_{2} p_{1}} \\
& =\frac{p_{1}^{2}\left(E-E_{1}\right)+\left(p_{1}^{2}-\boldsymbol{P} \cdot \boldsymbol{p}_{1}\right) E_{1}}{E_{1} E_{2} p_{1}} \\
& =\frac{p_{1}^{2} E-\left(\boldsymbol{P} \cdot \boldsymbol{p}_{1}\right) E_{1}}{E_{1} E_{2} p_{1}} \tag{B.7}
\end{align*}
$$

Substituting (B.7) in (B.3), we obtain the two-particle density of states.

$$
\begin{equation*}
\rho\left(E ; p_{1}, p_{2}\right)=\frac{E_{1} E_{2} p_{1}^{2}}{(2 \pi)^{3}\left\{E p_{1}^{2}-E_{1}\left(\boldsymbol{P} \cdot \boldsymbol{p}_{1}\right)\right\}} \tag{B.8}
\end{equation*}
$$

## Limiting case

If we allow the mass of the particle 2 to become infinite,

$$
m_{2} \rightarrow \infty, \text { then } E \rightarrow \infty \text { such that } \frac{E}{E_{2}} \rightarrow 1 ; \frac{E_{1}}{E_{2}} \rightarrow 0 .
$$

In this limiting case, the two-particle density of states (B.8)

$$
\begin{equation*}
\rho\left(E ; p_{1}, p_{2}\right) \longrightarrow \frac{E_{1} p_{1}}{(2 \pi)^{3}}, \tag{B.9}
\end{equation*}
$$

reduces to the one-particle density of states (B.1).

## 3. Three-particle final state

Here we have three particles in the final state with total energy $E$ and total momentum $\boldsymbol{P}$ such that

$$
E=E_{1}+E_{2}+E_{3} ; \quad \boldsymbol{P}=\boldsymbol{p}_{1}+\boldsymbol{p}_{2}+\boldsymbol{p}_{3} .
$$

The three-particle density of states is a product of one-particle density of states and the other two-particle density of states.

$$
\begin{align*}
\rho\left(E ; p_{1}, p_{2}, p_{3}\right) & =\rho\left(E_{1}, p_{1}\right) \rho\left(E-E_{1} ; p_{2}, p_{3}\right) \\
& =\frac{E_{1} p_{1}}{(2 \pi)^{3}} \frac{E_{2} E_{3} p_{2}^{3}}{(2 \pi)^{3}\left\{\left(E-E_{1}\right) p_{2}^{2}-E_{2}\left(\boldsymbol{P}-\boldsymbol{p}_{1}\right) \cdot \boldsymbol{p}_{2}\right\}} \\
& =\frac{1}{(2 \pi)^{6}} \frac{E_{1} E_{2} E_{3} p_{1} p_{2}^{3}}{\left(E-E_{1}\right) p_{2}^{2}-E_{2}\left(\boldsymbol{P}-\boldsymbol{p}_{1}\right) \cdot \boldsymbol{p}_{2}} . \tag{B.10}
\end{align*}
$$

## Limiting case

If $m_{3} \rightarrow \infty$, then $E_{3} \rightarrow \infty$. Consequently,

$$
\begin{equation*}
\rho\left(E ; p_{1}, p_{2}, p_{3}\right)=\frac{1}{(2 \pi)^{6}} E_{1} E_{2} p_{1} p_{2}, \tag{B.11}
\end{equation*}
$$

since $\frac{E-E_{1}}{E_{3}} \rightarrow 1$ and $\frac{E_{2}}{E_{3}} \rightarrow 0$. The three-particle density of states reduces to a product of two one-particle density of states when one of the particles has an infinite mass.

## 4. n -particle final state

The construction of three-particle density of states has given us a clue how to construct n-particle density of states. If $E$ is the total energy and $\boldsymbol{P}$, the total momentum of the n -particle state, then the n -particle density of states is the product of one-particle density of states with energy $E_{1}$ and momentum $p_{1}$ and the ( $\mathrm{n}-1$ )-particle density of states with energy $E-E_{1}$ and momentum $\boldsymbol{P}-\boldsymbol{p}_{1}$.

$$
\begin{equation*}
\rho\left(E ; p_{1}, p_{2}, \cdots, p_{n}\right)=\rho\left(E_{1}, p_{1}\right) \rho\left(E-E_{1} ; p_{2}, p_{3}, \cdots, p_{n}\right) . \tag{B.12}
\end{equation*}
$$

Thus, the n-particle density of states can be constructed by applying the rule successively.

For additional reading and for a comprehensive account on density of states, the reader is referred to: Alladi Ramakrishnan, Handbuch der Physik, Vol. III, p.524, Springer-Verlag (1959); Alladi Ramakrishnan, Elementary Particles and Cosmic Rays, p.503, Pergamon Press, Oxford (1962); R.P. Feynman, Quantum Electrodynamics, p.81, W.A. Benjamin, New York (1962).

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[^0]:    ${ }^{1}$ For further details, please refer to V. Devanathan, Quantum Mechanics, Narosa Publishing House, New Delhi (2005).

[^1]:    ${ }^{2}$ W. Pauli and V. Weisskopf, Helv. Phys. Acta. 7, 709 (1934).

[^2]:    ${ }^{3}$ It is the most commonly used representation but it is not the unique representation. Other representations are possible.

[^3]:    ${ }^{4}$ The space-time dependence of the Dirac wave function $\Psi$ is given by Eq. (1.24) and it does not affect $\rho$ and $\boldsymbol{J}$ as shown in Eq. (1.35).

[^4]:    ${ }^{5}$ Normally, we think of the vacuum as a state of emptiness but the Dirac vacuum has an infinite number of particles occupying the negative energy states!
    ${ }^{6}$ C. D. Anderson, Phys. Rev., 41, 405 (1932). C. D. Anderson was awarded the Nobel prize for physics in 1936 for the discovery of positron.

[^5]:    ${ }^{7}$ R. P. Feynman, Phys. Rev., 76, 749 (1949).

[^6]:    ${ }^{8}$ F. J. Dyson, Phys. Rev., 75, 486 (1949).

[^7]:    ${ }^{9}$ The matrices $\alpha_{x}, \alpha_{y}$ and $\alpha_{z}$ are independent of the position coordinates and momenta and hence commute with them.

[^8]:    ${ }^{10}$ In deducing Eq. (1.60), problem (11.6) of my book on Quantum Mechanics will be helpful.

[^9]:    ${ }^{1}$ In non-relativistic physics, we usually write $\psi^{\dagger} \mathcal{O} \psi$ but in relativistic physics, we should write $\bar{\psi} \mathcal{O} \psi$ since it is $\bar{\psi} \psi$ and not $\psi^{\dagger} \psi$, that transforms as a Lorentz scalar.

[^10]:    ${ }^{2}$ R. P. Feynman, Quantum Electrodynamics, W. A. Benjamin, Inc., New York (1962).
    ${ }^{3}$ Schweber et al. (1956) choose a slightly different normalization $\bar{\psi}_{p} \psi_{p}=1$ and $\bar{\psi}_{n} \psi_{n}=-1$ such that

    $$
    \sum_{\text {spins }}\left(\tilde{\psi}_{p} \psi_{p}-\tilde{\psi}_{n} \psi_{n}\right)=4 \quad \text { and } \sum_{\text {spins }}\left(\psi_{p} \tilde{\psi}_{p}-\psi_{n} \tilde{\psi}_{n}\right)=I
    $$

[^11]:    ${ }^{4}$ The two spin orientations, spin up and spin down are denoted by arrow marks $\uparrow$ and $\downarrow$ in Eqs. (2.59) and (2.60).

[^12]:    ${ }^{5}$ Feynman's definition of $\gamma_{5}=\gamma_{x} \gamma_{y} \gamma_{z} \gamma_{0}$ differs by a change in sign.

[^13]:    ${ }^{6}$ We use bold upright letters $\mathbf{a}, \mathbf{b}$ to denote the four-vectors to distinguish them from the three-vectors for which we use bold italics $\boldsymbol{a}, \boldsymbol{b}$.

[^14]:    ${ }^{1}$ The Dirac equation for the electron with negative charge $(-e)$ is

    $$
    \left\{c \boldsymbol{\alpha} \cdot\left(\boldsymbol{p}+\frac{e \boldsymbol{A}}{c}\right)+\beta m c^{2}\right\} \psi=(E+e \phi) \psi
    $$

[^15]:    ${ }^{2}$ It can be easily checked that if the magnetic field $\boldsymbol{B}$ is constant, then the relation

[^16]:    ${ }^{3}$ Starting from the relativistic relation between energy and momentum, we get

    $$
    \begin{aligned}
    E & =\left(p^{2} c^{2}+m^{2} c^{4}\right)^{1 / 2}=m c^{2}\left(1+\frac{p^{2}}{m^{2} c^{2}}\right)^{1 / 2} \\
    & =m c^{2}\left\{1+\frac{p^{2}}{2 m^{2} c^{2}}-\frac{p^{4}}{8 m^{4} c^{4}}+\cdots\right\}=m c^{2}+\frac{p^{2}}{2 m}-\frac{p^{4}}{8 m^{3} c^{2}}+\cdots
    \end{aligned}
    $$

[^17]:    ${ }^{4}$ The Hamiltonian $H$ consists of three terms, of which the terms involving $\boldsymbol{\alpha}$ and $\beta$ are $4 \times 4$ matrices but the potential term is a scalar and so it should be multiplied by $4 \times 4$ unit matrix. In general, in a matrix equation, all the terms should be treated as matrices of the same dimension.

[^18]:    ${ }^{5}$ An alternative form for the Dirac Hamiltonian in spherical polar coordinates is (vide Solved Problems 2.5 and 2.6.)

    $$
    H=\alpha_{r} p_{r}+\frac{i}{r} \alpha_{r} \beta K+\beta m+V(r)
    $$

    where $\alpha_{r}$ and $p_{r}$ are radial components of $\boldsymbol{\alpha}$ and $\boldsymbol{p}$, defined by

    $$
    \alpha_{r}=\frac{1}{r}(\boldsymbol{\alpha} \cdot \boldsymbol{r}), \quad \quad p_{r}=\frac{1}{r}(\boldsymbol{r} \cdot \boldsymbol{p}-i)
    $$

[^19]:    ${ }^{6}$ For a comprehensive discussion on Angular Momentum, the reader is referred to "Angular Momentum Techniques in Quantum Mechanics" by V. Devanathan, Kluwer Academic Publishers (1999).

[^20]:    ${ }^{7}$ W. E. Lamb, Jr and R. C. Retherford, Phys. Rev., 72, 241(1957).

[^21]:    ${ }^{8}$ V. Devanathan, Angular Momentum Techniques in Quantum Mechanics, Kluwer Academic Publishers (1999).

[^22]:    ${ }^{9}$ V. Devanathan, Nuclear Physics, Appendices C \& D, Narosa Publishing House, New Delhi (2006).

[^23]:    ${ }^{1}$ F. Reines and C.L. Cowan,Jr.," Measurement of free antineutrino absorption cross section by protons", Phys. Rev. 113, 273 (1959). F. Reines, "Neutrino Interactions", Ann. Rev. Nucl. Sci. 10, 1 (1960).
    ${ }^{2}$ M. Perl et al, Phys. Rev. Lett. 35, 1489 (1975).

[^24]:    ${ }^{3}$ E.C.G. Sudarshan and R.E. Marshak, Phys. Rev. 109, 860 (1958).
    ${ }^{4}$ R.P. Feynman and M. Gell-Mann, Phys. Rev. 109, 193 (1958).

[^25]:    ${ }^{5}$ For convenience of writing, we shall hereafter omit the suffix $W$ for the matrices with the implicit understanding that we are working in the Weyl representation.

[^26]:    ${ }^{6}$ J. Schechter and J.W.F. Valle, Phys. Rev. D25, 295 (1982).

[^27]:    ${ }^{7}$ S.R. Elliot, A.A. Hahn and M. Moe, Phys. Rev. Lett. 59, 2020 (1987).
    ${ }^{8}$ H.L. Harney, H.V. Klapdor-Kleingrothans, A. Dietz and I.V. Krivosheina, Mod. Phys. Lett. A16, 2409 (2001); H.V. Klapdor-Kleingrothans and I.V. Krivosheina, Mod. Phys. Lett. A21, 1547 (2006).

[^28]:    ${ }^{9}$ J.N. Bahcall, Neutrino Astrophysics, Cambridge University Press (1989).

[^29]:    ${ }^{10}$ R. Davis, D.S. Harmer and K.C. Hoffman, Phys. Rev. Lett. 20, 1205 (1968).
    ${ }^{11}$ Kamiokande Collaboration: K.S. Hirata et al, Phys. Rev. Lett. 63, 16 (1989); Phys. Lett. B280, 146 (1992).

[^30]:    ${ }^{12}$ SNO Collaboration: Q.R. Ahmad et al, Phys. Rev. Lett. 87, 071301 (2001).

[^31]:    ${ }^{13}$ Super-Kamiokande Collaboration: Y. Fukuda et al, Phys. Rev. Lett. 81, 1562 (1998).

[^32]:    ${ }^{14}$ R. Parthasarathy, Relativistic Quantum Mechanics, Narosa Publishing House, New Delhi (2010).

[^33]:    ${ }^{15}$ Distance of the earth from Sun at perihelion $=1.47 \times 10^{11} \mathrm{~m}$.
    Distance of the earth from Sun at aphelion $=1.52 \times 10^{11} \mathrm{~m}$.
    Mean distance of the earth from Sun $=1.497 \times 10^{11} \mathrm{~m}$.
    ${ }^{16}$ M. Narayan, M.V.N. Murthy, G. Rajasekaran and S. Uma Sankar, Phys. Rev., D53, 2809 (1996); M. Narayan, G. Rajasekaran and S. Uma Sankar, Phys. Rev., D56, 437 (1997); D. Indhumathi, M.V.N. Murthy, G. Rajasekaran, Nita Sinha, Phys. Rev., D74, 53004 (2006).

[^34]:    ${ }^{17}$ Kobayashi and Maskawa has suggested that the dimension of the unitary matrix has to be atleast 3 for a phase to exist. It is the CP violating phase which signals the matter-antimatter asymmetric universe that is of great cosmological importance. Kobayashi and Maskawa were awarded the Nobel Prize in the year 2008.

[^35]:    ${ }^{1}$ V. Devanathan, Quantum Mechanics, (Ch.9, Theory of Scattering) Narosa Publishers, New Delhi (2005)

[^36]:    ${ }^{2}$ S. S. Schweber and H. A. Bethe, Mesons and Fields, Vol. 1, Row, Peterson \& Co., New York (1956).

[^37]:    ${ }^{3}$ For brevity, we use the notation 1 for $\left(\boldsymbol{x}_{1}, t_{1}\right)$ and $\mathbf{2}$ for $\left(\boldsymbol{x}_{2}, t_{2}\right)$.
    ${ }^{4}$ The suffixes $1,2,3$ are used to denote the space-time coordinates and also the $\mathrm{x}, \mathrm{y}, \mathrm{z}$ components but the reader can easily infer the correct meaning from the context they are used.
    ${ }^{5}$ The Dirac equation for a free particle in natural units can be written either as

    $$
    (p-m) \Psi=0 \quad \text { or } \quad(i \not \nabla-m) \Psi=0
    $$

[^38]:    ${ }^{6}$ We use bold italics $\boldsymbol{p}, \boldsymbol{x}$ to denote three-dimensional vectors and upright bold letters $\mathbf{p}, \mathbf{x}$ to denote four vectors.

[^39]:    ${ }^{1}$ We denote the four-vectors by upright bold letters ( $\mathbf{p}, \mathbf{x}$ ), the three-vectors by italic bold letters $(\boldsymbol{p}, \boldsymbol{x})$ and the scalars by ordinary italics $(p, x)$.

[^40]:    ${ }^{2}$ The convention of summation over repeated indices is used. Greek letters $\mu, \nu, \lambda, \rho$ are used as suffixes to denote the components of four-vectors and Latin letters $i, j, k, l$ are used as suffixes to denote the components of three-vectors.

[^41]:    ${ }^{3}$ The integral representation of the delta function $\delta(x)$ is used.

    $$
    \delta(x)=\frac{1}{2 \pi} \int_{-\infty}^{+\infty} e^{i x t} d t
    $$

    ${ }^{4}$ The reader is advised to refer to Sec. 2.2 (Algebra of $\gamma$ matrices) for details of evaluation of traces.

[^42]:    ${ }^{5}$ Underbraces are used to highlight the factors that are considered in the step-by-step calculation for the benefit of the reader. Note that $\gamma_{0} p_{i}=-p_{i} \gamma_{0}+2 E_{i}$.
    ${ }^{6} E^{2}=p^{2}+m^{2}$. Differentiating, we get $2 E d E=2 p d p$. Therefore $v=d E / d p=p / E$.

[^43]:    ${ }^{7}$ The reader is referred to Appendix B for a comprehensive discussion on density of final states.
    ${ }^{8}$ O. Klein and Y. Nishina, Z. Physik, 52, 853 (1929).

[^44]:    ${ }^{9}$ In Eq. (6.51),$-\mathcal{O}$ is replaced by $\mathcal{O}$, since the negative sign is of no consequence in the calculation of transition probability and hence omitted.

[^45]:    ${ }^{10}$ In finding the absolute square of the matrix element, the summation index $\mu$ has to be changed to $\nu$ in the hermitian conjugate of the matrix element, since

    $$
    \left|\sum_{\mu} a_{\mu}\right|^{2}=\left(\sum_{\mu} a_{\mu}\right)\left(\sum_{\nu} a_{\nu}^{*}\right)=\sum_{\mu \nu} a_{\mu} a_{\nu}^{*}
    $$

[^46]:    ${ }^{11}$ A summation over the repeated index is implied.
    ${ }^{12}$ Underbraces are used to highlight the factors that are considered in the step-by-step calculation for the benefit of the reader.

[^47]:    ${ }^{13}$ C. Mфller, Ann. der. Phys., 85, 711 (1932).
    ${ }^{14}$ H. J. Bhabha, Proc. Roy. Soc (London)A154, 195 (1936).

[^48]:    ${ }^{15}$ R. P. Feynman, Quantum Electrodynamics W. A. Benjamin (1976).

[^49]:    ${ }^{16} \mathrm{~A}$ summation over the repeated index is implied.

[^50]:    ${ }^{17}$ Underbraces are used to highlight the terms that are used in the step-by-step calculation.

[^51]:    ${ }^{18}$ We use bold letters $\mathbf{a}, \mathbf{b}, \mathbf{c}, \mathbf{d}$ to denote the four-vectors and bold italic letters $\boldsymbol{a}, \boldsymbol{b}, \boldsymbol{c}, \boldsymbol{d}$ to denote the three vectors.

[^52]:    ${ }^{19}$ We have used underbrace to highlight the relevant terms in the step-by-step evaluation.

[^53]:    ${ }^{20} \varepsilon_{1}$ has no time component and its space component is normal to the momentum vector $\boldsymbol{k}_{1}$.

[^54]:    ${ }^{1}$ R. P. Feynman, Quantum Electrodynamics, W. A. Benjamin, Inc. (1962).
    ${ }^{2}$ We use upright bold letters ( $\mathbf{k}, \mathbf{p}$ ) to denote four-vectors, italic bold letters ( $\boldsymbol{k}, \boldsymbol{p}$ ) to denote three-vectors. The symbols $\not p, \nless k$ represent $\not p=\gamma_{\mu} p_{\mu}, \not \not k=\gamma_{\mu} k_{\mu}$.

[^55]:    ${ }^{3}$ The theoretical mass $m_{\mathrm{th}}$ is the bare mass $m_{0}$ without self-interaction and the experimental mass $m_{\exp }$ is the observed mass $m$ of the electron.

[^56]:    ${ }^{1}$ For greater details, the reader is referred to: H. Goldstein, Classical Mechanics, Narosa Publishing House, New Delhi (Indian Edition) (2001).
    ${ }^{2}$ Of the various paths available between the two given end points, the system chooses the path for which the action integral $\mathscr{A}$ is a minimum. This is found by the method of calculus of variations.

[^57]:    ${ }^{3}$ The Minkowski space is a natural extension of the ordinary three dimensional Cartesian space with which we are familiar, to the four dimensional space-time which is required for any theory that includes the special theory of relativity. This avoids the use of metric tensor and the artificial distinction between the covariant and contravariant quantities. However this is not possible for the treatment of the general theory of relativity with the concept of curved space. So, many authors prefer to use the metric $1,-1,-1,-1$ with the accompanying distinction between covariant and contravariant quantities.

[^58]:    ${ }^{4}$ The square brackets are used to denote the commutation relation and the curly brackets denote the anticommutation relation.

    $$
    \left[a_{k}, a_{l}^{\dagger}\right]_{-}=a_{k} a_{l}^{\dagger}-a_{l}^{\dagger} a_{k}=\delta_{k l} ; \quad\left\{a_{k}, a_{l}^{\dagger}\right\}_{+}=a_{k} a_{l}^{\dagger}+a_{l}^{\dagger} a_{k}=\delta_{k l}
    $$

[^59]:    ${ }^{5}$ P. Jordan and E. P. Wigner, Z. Physik, 47, 631 (1928).

[^60]:    ${ }^{1}$ W. Pauli and V. Weisskopf, Helv. Phys. Acta., 7, 709 (1934).

[^61]:    ${ }^{2}$ The expansion coefficients are denoted either by $a_{k}, a_{k}^{\dagger}$ or by $a(\boldsymbol{k}), a^{\dagger}(\boldsymbol{k})$.

[^62]:    ${ }^{3}$ Since we are considering only the K-G field in his chapter, we will hereafter omit the suffix $K G$ for the $\mathscr{L}$.

[^63]:    ${ }^{4}$ The reader is referred to: V. Devanathan, Quantum Mechanics, Sec. 6.4.3, Narosa Publishing House, New Delhi (2005), for further details on simple harmonic oscillator.

[^64]:    ${ }^{3}$ Since $\psi \bar{\psi}$ is a matrix and $\bar{\psi} \psi$ is a number, there is no meaning in trying to find the anticommutation relation $\{\psi(\boldsymbol{x}), \bar{\psi}(\boldsymbol{x})\}_{+}$. We can only find the anticommutation relation between their components $\left\{\psi_{\alpha}(\boldsymbol{x}), \bar{\psi}_{\beta}(\boldsymbol{x})\right\}_{+}$.

[^65]:    ${ }^{2}$ The symbol $\epsilon_{i j k}$ stands for
    $\epsilon_{i j k}=+1, \quad$ if $i, j, k$ is an even permutation of the numbers $1,2,3$
    $=-1, \quad$ if $i, j, k$ is an odd permutation of the numbers $1,2,3$
    $=0, \quad$ if any two indices are equal.

[^66]:    ${ }^{3}$ The electromagnetic field in a space region which is far removed from any sources of the field is also called the radiation field.
    ${ }^{4}$ S.N. Gupta, Proc. Phys. Soc. A63, 681 (1950); K. Bleuler, Helv. Phys. Acta 23, 567 (1950).

[^67]:    ${ }^{5}$ The four-vectors are denoted by upright bold face letters like $\mathbf{k}, \mathbf{x}$ and the three vectors by italic bold face letters like $\boldsymbol{k}, \boldsymbol{x}$.

[^68]:    ${ }^{6}$ The annihilation and creation operators are represented either as $a^{(\lambda)}, a^{(\lambda) \dagger}$ or as $a_{\lambda}, a_{\lambda}^{\dagger}$, using $\lambda$ as a superscript or as a subscript, depending upon convenience.

[^69]:    ${ }^{1}$ The reader may refer to chap. 6 and 7 of my book: V. Devanathan, Quantum Mechanics, Narosa Publishing House, New Delhi (2005), for a better understanding.

[^70]:    ${ }^{2}$ S.S. Schweber, H. A. Bethe and F. D. Hoffmann, Mesons and Fields, Vol 1, p.167169, Row, Peterson and Co. (1956).

[^71]:    ${ }^{3}$ F. J. Dyson, Phys. Rev., 75, 486, 1736 (1949).

[^72]:    ${ }^{4}$ F. J. Dyson, Phys. Rev., 82, 428 (1951); 83, 608 (1951).
    ${ }^{5}$ G.C. Wick, Phys. Rev., 80, 268 (1950).

[^73]:    ${ }^{6}$ It may be recalled that the positive frequency operators $\phi^{+}, \psi^{+}, \bar{\psi}^{+}$consist of annihilation operators whereas the negative frequency operators $\phi^{-}, \psi^{-}, \bar{\psi}^{-}$carry creation operators.

[^74]:    ${ }^{1}$ Since the electromagnetic field strength tensor $F_{\mu \nu}$

    $$
    F_{\mu \nu}=\partial_{\mu} A_{\nu}-\partial_{\nu} A_{\mu}
    $$

[^75]:    ${ }^{2}$ The gauge covariant derivative of $\psi, D_{\mu}$ transforms in the same way as $\psi$.

    $$
    \psi(\mathbf{x}) \rightarrow e^{i \alpha(\mathbf{x})} \psi(\mathbf{x}) ; \quad D_{\mu} \psi(\mathbf{x}) \rightarrow e^{i \alpha(\mathbf{x})} D_{\mu} \psi(\mathbf{x})
    $$

[^76]:    ${ }^{3}$ For a good account on this subject, the reader is referred to: F. Halzen and A. D. Martin, Quarks and Leptons, John Wiley (1984).

[^77]:    ${ }^{4}$ F. Englert and R. Brout, Phys.Rev. Lett. 13, 321-323 (1964).
    ${ }^{5}$ Peter W. Higgs, Phys. Rev. Lett. 13, 508-509 (1964).
    ${ }^{6}$ G.S. Guralink, C.R. Hagen and T.W.B. Kibbs, Phys. Rev. Lett. 13, 585-587 (1964).

[^78]:    ${ }^{7}$ Although the standard model of elementary particles assumes zero mass for the neutrinos, there is a strong evidence that the neutrinos should have a small mass due to the discovery of neutrino oscillations.

[^79]:    ${ }^{8}$ t Hooft, G. (1971). "Renormalization of massless Yang-Mills fields". Nuclear Physics B 33, 173-177 (1971); "Renormalizable Lagrangians for massive Yang-Mills fields". Nuclear Physics B 35 167-448 (1971).

[^80]:    ${ }^{1}$ The integral $\int^{\infty} d r r^{n}$ diverges linearly for $n=0$, quadratically for $n=1$ and so on and $\int^{\infty} \frac{d r}{r}$ diverges logarithmically.

[^81]:    ${ }^{2}$ Consider a guitar string that has been tuned by stretching the string under tension. Depending on how the string is plucked and how much tension is in the string, different musical notes will be created by the string. These musical notes could be said to be excitation modes of that guitar string under tension.

    In a similar manner, in string theory, the elementary particles we observe could be thought of as the "musical notes" or excitation modes of elementary strings.
    ${ }^{3}$ J. Scherk and J. H. Schwarz, Nucl. Phys. B81, 118 (1974).

[^82]:    ${ }^{4}$ A. Ramadevi, A Pedagogical Introduction to String Theory, Physics News, 40, 19 (2010)

